# INVARIANCE THEORY, THE HEAT EQUATION, 

 AND THE
## ATIYAH-SINGER INDEX THEOREM

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## INTRODUCTION

This book treats the Atiyah-Singer index theorem using heat equation methods. The heat equation gives a local formula for the index of any elliptic complex. We use invariance theory to identify the integrand of the index theorem for the four classical elliptic complexes with the invariants of the heat equation. Since the twisted signature complex provides a sufficiently rich family of examples, this approach yields a proof of the Atiyah-Singer theorem in complete generality. We also use heat equation methods to discuss Lefschetz fixed point formulas, the Gauss-Bonnet theorem for a manifold with smooth boundary, and the twisted eta invariant. We shall not include a discussion of the signature theorem for manifolds with boundary.

The first chapter reviews results from analysis. Sections 1.1 through 1.7 represent standard elliptic material. Sections 1.8 through 1.10 contain the material necessary to discuss Lefschetz fixed point formulas and other topics.

Invariance theory and differential geometry provide the necessary link between the analytic formulation of the index theorem given by heat equation methods and the topological formulation of the index theorem contained in the Atiyah-Singer theorem. Sections 2.1 through 2.3 are a review of characteristic classes from the point of view of differential forms. Section 2.4 gives an invariant-theoretic characterization of the Euler form which is used to give a heat equation proof of the Gauss-Bonnet theorem. Sections 2.5 and 2.6 discuss the Pontrjagin forms of the tangent bundle and the Chern forms of the coefficient bundle using invariance theory.

The third chapter combines the results of the first two chapters to prove the Atiyah-Singer theorem for the four classical elliptic complexes. We first present a heat equation proof of the Hirzebruch signature theorem. The twisted spin complex provides a unified way of discussing the signature, Dolbeault, and de Rham complexes. In sections 3.2-3.4, we discuss the half-spin representations, the spin complex, and derive a formula for the $\hat{A}$ genus. We then discuss the Riemann-Roch formula for an almost complex manifold in section 3.5 using the SPIN $_{c}$ complex. In sections $3.6-3.7$ we give a second derivation of the Riemann-Roch formula for holomorphic Kaehler manifods using a more direct approach. In the final two sections we derive the Atiyah-Singer theorem in its full generality.

The final chapter is devoted to more specialized topics. Sections 4.1-4.2 deal with elliptic boundary value problems and derive the Gauss-Bonnet theorem for manifolds with boundary. In sections 4.3-4.4 we discuss the twisted eta invariant on a manifold without boundary and we derive the Atiyah-Patodi-Singer twisted index formula. Section 4.5 gives a brief discussion of Lefschetz fixed point formulas using heat equation methods. In section 4.6 we use the eta invariant to calculate the $K$-theory of spherical space forms. In section 4.7, we discuss Singer's conjecture for the Euler form and related questions. In section 4.8 , we discuss the local formulas for the invariants of the heat equation which have been derived by several authors, and in section 4.9 we apply these results to questions of spectral geometry.

The bibliography at the end of this book is not intended to be exhaustive but rather to provide the reader with a list of a few of the basic papers which have appeared. We refer the reader to the bibliography of Berger and Berard for a more complete list of works on spectral geometry.

This book is organized into four chapters. Each chapter is divided into a number of sections. Each Lemma or Theorem is indexed according to this subdivision. Thus, for example, Lemma 1.2.3 is the third Lemma of section 2 of Chapter 1.

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## CHAPTER 1 <br> PSEUDO-DIFFERENTIAL OPERATORS

## Introduction

In the first chapter, we develop the analysis needed to define the index of an elliptic operator and to compute the index using heat equation methods. Sections 1.1 and 1.2 are brief reviews of Sobolev spaces and pseudodifferential operators on Euclidean spaces. In section 1.3, we transfer these notions to compact Riemannian manifolds using partition of unity arguments. In section 1.4 we review the facts concerning Fredholm operators needed in section 1.5 to prove the Hodge decomposition theorem and to discuss the spectral theory of self-adjoint elliptic operators. In section 1.6 we introduce the heat equation and in section 1.7 we derive the local formula for the index of an elliptic operator using heat equation methods. Section 1.8 generalizes the results of section 1.7 to find a local formula for the Lefschetz number of an elliptic complex. In section 1.9, we discuss the index of an elliptic operator on a manifold with boundary and in section 1.10, we discuss the zeta and eta invariants.

Sections 1.1 and 1.4 review basic facts we need, whereas sections 1.8 through 1.10 treat advanced topics which may be omitted from a first reading. We have attempted to keep this chapter self-contained and to assume nothing beyond a first course in analysis. An exception is the de Rham theorem in section 1.5 which is used as an example.

A number of people have contributed to the mathematical ideas which are contained in the first chapter. We were introduced to the analysis of sections 1.1 through 1.7 by a course taught by L. Nirenberg. Much of the organization in these sections is modeled on his course. The idea of using the heat equation or the zeta function to compute the index of an elliptic operator seems to be due to R. Bott. The functional calculus used in the study of the heat equation contained in section 1.7 is due to R. Seeley as are the analytic facts on the zeta and eta functions of section 1.10.

The approach to Lefschetz fixed point theorems contained in section 1.8 is due to T. Kotake for the case of isolated fixed points and to S. C. Lee and the author in the general case. The analytic facts for boundary value problems discussed in section 1.9 are due to P. Greiner and R. Seeley.

### 1.1. Fourier Transform, Schwartz Class, And Sobolev Spaces.

The Sobolev spaces and Fourier transform provide the basic tools we shall need in our study of elliptic partial differential operators. Let $x=$ $\left(x_{1}, \ldots, x_{m}\right) \in \mathbf{R}^{m}$. If $x, y \in \mathbf{R}^{m}$, we define:

$$
x \cdot y=x_{1} y_{1}+\cdots+x_{m} y_{m} \quad \text { and } \quad|x|=(x \cdot x)^{1 / 2}
$$

as the Euclicean dot product and length. Let $\alpha=\left(\alpha_{1}, \ldots, \alpha_{m}\right)$ be a multiindex. The $\alpha_{j}$ are non-negative integers. We define:

$$
|\alpha|=\alpha_{1}+\cdots+\alpha_{m}, \quad \alpha!=\alpha_{1}!\ldots \alpha_{m}!, \quad x^{\alpha}=x_{1}^{\alpha_{1}} \ldots x_{m}^{\alpha_{m}}
$$

Finally, we define:

$$
d_{x}^{\alpha}=\left(\frac{\partial}{\partial x_{1}}\right)^{\alpha_{1}} \cdots\left(\frac{\partial}{\partial x_{m}}\right)^{\alpha_{m}} \quad \text { and } \quad D_{x}^{\alpha}=(-i)^{|\alpha|} d_{x}^{\alpha}
$$

as a convenient notation for multiple partial differentiation. The extra factors of $(-i)$ defining $D_{x}^{\alpha}$ are present to simplify later formulas. If $f(x)$ is a smooth complex valued function, then Taylor's theorem takes the form:

$$
f(x)=\sum_{|\alpha| \leq n} d_{x}^{\alpha} f\left(x_{0}\right) \frac{\left(x-x_{0}\right)^{\alpha}}{\alpha!}+O\left(\left|x-x_{0}\right|^{n+1}\right)
$$

The Schwartz class $\mathcal{S}$ is the set of all smooth complex valued functions $f$ on $\mathbf{R}^{m}$ such that for all $\alpha, \beta$ there are constants $C_{\alpha, \beta}$ such that

$$
\left|x^{\alpha} D_{x}^{\beta} f\right| \leq C_{\alpha, \beta} .
$$

This is equivalent to assuming there exist estimates of the form:

$$
\left|D_{x}^{\beta} f\right| \leq C_{n, \beta}(1+|x|)^{-n}
$$

for all $(n, \beta)$. The functions in $\mathcal{S}$ have all their derivatives decreasing faster at $\infty$ than the inverse of any polynomial.

For the remainder of Chapter 1, we let $d x, d y, d \xi$, etc., denote Lebesgue measure on $\mathbf{R}^{m}$ with an additional normalizing factor of $(2 \pi)^{-m / 2}$. With this normalization, the integral of the Gaussian distribution becomes:

$$
\int e^{-\frac{1}{2}|x|^{2}} d x=1
$$

We absorb the normalizing constant into the measure in order to simplify the formulas of the Fourier transform. If $C_{0}^{\infty}\left(\mathbf{R}^{m}\right)$ denotes the set of smooth functions of compact support on $\mathbf{R}^{m}$, then this is a subset of $\mathcal{S}$. Since $C_{0}\left(\mathbf{R}^{m}\right)$ is dense in $L^{2}\left(\mathbf{R}^{m}\right), \mathcal{S}$ is dense in $L^{2}\left(\mathbf{R}^{m}\right)$.

We define the convolution product of two elements of $\mathcal{S}$ by:

$$
(f * g)(x)=\int f(x-y) g(y) d y=\int f(y) g(x-y) d y
$$

This defines an associative and commutative multiplication. Although there is no identity, there do exist approximate identities:

Lemma 1.1.1. Let $f \in \mathcal{S}$ with $\int f(x) d x=1$. Define $f_{u}(x)=u^{-m} f\left(\frac{x}{u}\right)$. Then for any $g \in \mathcal{S}, f_{u} * g$ converges uniformly to $g$ as $u \rightarrow 0$.
Proof: Choose $C$ so $\int|f(x)| d x \leq C$ and $|g(x)| \leq C$. Because the first derivatives of $g$ are uniformly bounded, $g$ is uniformly continuous. Let $\varepsilon>0$ and choose $\delta>0$ so $|x-y| \leq \delta$ implies $|g(x)-g(y)| \leq \varepsilon$. Because $\int f_{u}(x) d x=1$, we compute:

$$
\begin{aligned}
\left|f_{u} * g(x)-g(x)\right| & =\left|\int f_{u}(y)\{g(x-y)-g(x)\} d y\right| \\
& \leq \int\left|f_{u}(y)\{g(x-y)-g(x)\}\right| d y
\end{aligned}
$$

We decompose this integral into two pieces. If $|y| \leq \delta$ we bound it by $C \varepsilon$. The integral for $|y| \geq \delta$ can be bounded by:

$$
2 C \int_{|y| \geq \delta}\left|f_{u}(y)\right| d y=2 C \int_{|y| \geq \delta / u}|f(y)| d y .
$$

This converges to zero as $u \rightarrow 0$ so we can bound this by $C \varepsilon$ if $u<u(\varepsilon)$. This completes the proof.

A similar convolution smoothing can be applied to approximate any element of $L^{p}$ arbitrarily well in the $L^{p}$ norm by a smooth function of compact support.

We define the Fourier transform $\hat{f}(\xi)$ by:

$$
\hat{f}(\xi)=\int e^{-i x \cdot \xi} f(x) d x \quad \text { for } f \in \mathcal{S}
$$

For the moment $\xi \in \mathbf{R}^{m}$; when we consider operators on manifolds, it will be natural to regard $\xi$ as an element of the fiber of the cotangent space. By integrating by parts and using Lebesgue dominated convergence, we compute:

$$
D_{\xi}^{\alpha}\{\hat{f}(\xi)\}=(-1)^{|\alpha|}\left\{\widehat{x^{\alpha} f}\right\} \quad \text { and } \quad \xi^{\alpha} \hat{f}(\xi)=\left\{\widehat{D_{x}^{\alpha} f}\right\} .
$$

This implies $\hat{f} \in \mathcal{S}$ so Fourier transform defines a map $\mathcal{S} \rightarrow \mathcal{S}$.
We compute the Fourier transform of the Gaussian distribution. Let $f_{0}(x)=\exp \left(-\frac{1}{2}|x|^{2}\right)$, then $f_{0} \in \mathcal{S}$ and $\int f_{0}(x) d x=1$. We compute:

$$
\begin{aligned}
\hat{f}_{0}(\xi) & =\int e^{-i x \cdot \xi} e^{-\frac{1}{2}|x|^{2}} d x \\
& =e^{-\frac{1}{2}|\xi|^{2}} \int e^{-(x+i \xi) \cdot(x+i \xi) / 2} d x
\end{aligned}
$$

We make a change of variables to replace $x+i \xi$ by $x$ and to shift the contour in $\mathbf{C}^{m}$ back to the original contour $\mathbf{R}^{m}$. This shows the integral is 1 and $\hat{f}_{0}(\xi)=\exp \left(-\frac{1}{2}|\xi|^{2}\right)$ so the function $f_{0}$ is its own Fourier transform.

In fact, the Fourier transform is bijective and the Fourier inversion formula gives the inverse expressing $f$ in terms of $\hat{f}$ by:

$$
f(x)=\int e^{i x \cdot \xi} \hat{f}(\xi) d \xi=\hat{\hat{f}}(-x)
$$

We define $T(f)=\hat{\hat{f}}(-x)=\int e^{i x \cdot \xi} \hat{f}(\xi) d \xi$ as a linear map from $\mathcal{S} \rightarrow \mathcal{S}$. We must show that $T(f)=f$ to prove the Fourier inversion formula.

Suppose first $f(0)=0$. We expand:

$$
f(x)=\int_{0}^{1} \frac{d}{d t}\{f(t x)\} d t=\sum_{j} x_{j} \int_{0}^{1} \frac{\partial f}{\partial x_{j}}(t x) d t=\sum_{j} x_{j} g_{j}
$$

where the $g_{j}$ are smooth. Let $\phi \in C_{0}^{\infty}\left(\mathbf{R}^{m}\right)$ be identically 1 near $x=0$. Then we decompose:

$$
f(x)=\phi f(x)+(1-\phi) f(x)=\sum_{j} x_{j} \phi g_{j}+\sum_{j} x_{j}\left\{\frac{x_{j}(1-\phi) f}{|x|^{2}}\right\}
$$

Since $\phi g_{j}$ has compact support, it is in $\mathcal{S}$. Since $\phi$ is identically 1 near $x=0, \quad x_{j}(1-\phi) f /|x|^{2} \in \mathcal{S}$. Thus we can decompose $f=\sum x_{j} h_{j}$ for $h_{j} \in \mathcal{S}$. We Fourier transform this identity to conclude:

$$
\hat{f}=\sum_{j}\left\{\widehat{x_{j} h_{j}}\right\}=\sum_{j} i \frac{\partial \hat{h}_{j}}{\partial \xi_{j}}
$$

Since this is in divergence form, $T(f)(0)=\int \hat{f}(\xi) d \xi=0=f(0)$.
More generally, let $f \in \mathcal{S}$ be arbitrary. We decompose $f=f(0) f_{0}+$ $\left(f-f(0) f_{0}\right)$ for $f_{0}=\exp \left(-\frac{1}{2}|x|^{2}\right)$. Since $\hat{f}_{0}=f_{0}$ is an even function, $T\left(f_{0}\right)=f_{0}$ so that $T(f)(0)=f(0) f_{0}(0)+T\left(f-f(0) f_{0}\right)=f(0) f_{0}(0)=f(0)$ since $\left(f-f(0) f_{0}\right)(0)=0$. This shows $T(f)(0)=f(0)$ in general.

We use the linear structure on $\mathbf{R}^{m}$ to complete the proof of the Fourier inversion formula. Let $x_{0} \in \mathbf{R}^{m}$ be fixed. We let $g(x)=f\left(x+x_{0}\right)$ then:

$$
\begin{aligned}
f\left(x_{0}\right)=g(0)=T(g)(0) & =\int e^{-i x \cdot \xi} f\left(x+x_{0}\right) d x d \xi \\
& =\int e^{-i x \cdot \xi} e^{i x_{0} \cdot \xi} f(x) d x d \xi \\
& =T(f)\left(x_{0}\right)
\end{aligned}
$$

This shows the Fourier transform defines a bijective map $\mathcal{S} \rightarrow \mathcal{S}$. If we use the constants $C_{\alpha, \beta}=\sup _{x \in \mathbf{R}^{m}}\left|x^{\alpha} D_{x}^{\beta} f\right|$ to define a Frechet structure on $\mathcal{S}$, then the Fourier transform is a homeomorphism of topological vector spaces. It is not difficult to show $C_{0}^{\infty}\left(\mathbf{R}^{m}\right)$ is a dense subset of $\mathcal{S}$ in this topology. We can use either pointwise multiplication or convolution to define a multiplication on $\mathcal{S}$ and make $\mathcal{S}$ into a ring. The Fourier transform interchanges these two ring structures. We compute:

$$
\begin{aligned}
\hat{f} \cdot \hat{g} & =\int e^{-i x \cdot \xi} f(x) e^{-i y \cdot \xi} g(y) d x d y \\
& =\int e^{-i(x-y) \cdot \xi} f(x-y) e^{-i y \cdot \xi} g(y) d x d y \\
& =\int e^{-i x \cdot \xi} f(x-y) g(y) d x d y
\end{aligned}
$$

The integral is absolutely convergent so we may interchange the order of integration to compute $\hat{f} \cdot \hat{g}=(\widehat{f * g})$. If we replace $f$ by $\hat{f}$ and $g$ by $\hat{g}$ we see $(f \cdot g)(-x)=(\widehat{\hat{f} * \hat{g}})$ using the Fourier inversion formula. We now take the Fourier transform and use the Fourier inversion formula to see $(\widehat{f \cdot g})(-\xi)=(\hat{f} * \hat{g})(-\xi)$ so that $(\widehat{f \cdot g})=\hat{f} * \hat{g}$.

The final property we shall need of the Fourier transform is related to the $L^{2}$ inner product $(f, g)=\int f(x) \bar{g}(x) d x$. We compute:

$$
\begin{aligned}
(\hat{f}, g) & =\int f(x) e^{-i x \cdot \xi} \bar{g}(\xi) d x d \xi=\int f(x) e^{-i x \cdot \xi} \bar{g}(\xi) d \xi d x \\
& =(f, \hat{g}(-x))
\end{aligned}
$$

If we replace $g$ by $\hat{g}$ then $(\hat{f}, \hat{g})=(f, \hat{\hat{g}}(-x))=(f, g)$ so the Fourier transform is an isometry with respect to the $L^{2}$ inner product. Since $\mathcal{S}$ is dense in $L^{2}$, it extends to a unitary map $L^{2}\left(\mathbf{R}^{m}\right) \rightarrow L^{2}\left(\mathbf{R}^{m}\right)$. We summarize these properties of the Fourier transform as follows:
Lemma 1.1.2. The Fourier transform is a homeomorphism $\mathcal{S} \rightarrow \mathcal{S}$ such that:
(a) $f(x)=\int e^{i x \cdot \xi} \hat{f}(\xi) d \xi=\int e^{i(x-y) \cdot \xi} f(y) d y d \xi$ (Fourier inversion formula);
(b) $D_{x}^{\alpha} f(x)=\int e^{i x \cdot \xi} \xi^{\alpha} \hat{f}(\xi) d \xi$ and $\xi^{\alpha} \hat{f}(\xi)=\int e^{-i x \cdot \xi} D_{x}^{\alpha} f(x) d x$;
(c) $\hat{f} \cdot \hat{g}=(\widehat{f * g})$ and $\hat{f} * \hat{g}=(\widehat{f \cdot g})$;
(d) The Fourier transform extends to a unitary map of $L^{2}\left(\mathbf{R}^{m}\right) \rightarrow L^{2}\left(\mathbf{R}^{m}\right)$ such that $(f, g)=(\hat{f}, \hat{g})$. (Plancherel theorem).

We note that without the normalizing constant of $(2 \pi)^{-m / 2}$ in the definition of the measures $d x$ and $d \xi$ there would be various normalizing constants appearing in these identities. It is property (b) which will be of the
most interest to us since it will enable us to interchange differentiation and multiplication.

We define the Sobolev space $H_{s}\left(\mathbf{R}^{m}\right)$ to measure $L^{2}$ derivatives. If $s$ is a real number and $f \in \mathcal{S}$, we define:

$$
|f|_{s}^{2}=\int\left(1+|\xi|^{2}\right)^{s}|\hat{f}(\xi)|^{2} d \xi
$$

The Sobolev space $H_{s}\left(\mathbf{R}^{m}\right)$ is the completion of $\mathcal{S}$ with respect to the norm $\left.\right|_{s}$. The Plancherel theorem shows $H_{0}\left(\mathbf{R}^{m}\right)$ is isomorphic to $L^{2}\left(\mathbf{R}^{m}\right)$. More generally, $H_{s}\left(\mathbf{R}^{m}\right)$ is isomorphic to $L^{2}$ with the measure $\left(1+|\xi|^{2}\right)^{s / 2} d \xi$. Replacing $\left(1+|\xi|^{2}\right)^{s}$ by $(1+|\xi|)^{2 s}$ in the definition of $\left.\right|_{s}$ gives rise to an equivalent norm since there exist positive constants $c_{i}$ such that:

$$
c_{1}\left(1+|\xi|^{2}\right)^{s} \leq(1+|\xi|)^{2 s} \leq c_{2}\left(1+|\xi|^{2}\right)^{s} .
$$

In some sense, the subscript " $s$ " counts the number of $L^{2}$ derivatives. If $s=n$ is a positive integer, there exist positive constants $c_{1}, c_{2}$ so:

$$
c_{1}\left(1+|\xi|^{2}\right)^{n} \leq \sum_{|\alpha| \leq n}\left|\xi^{\alpha}\right|^{2} \leq c_{2}\left(1+|\xi|^{2}\right)^{n}
$$

This implies that we could define

$$
|f|_{n}^{2}=\sum_{|\alpha| \leq n} \int\left|\xi^{\alpha} \hat{f}\right|^{2} d \xi=\sum_{|\alpha| \leq n} \int\left|D_{x}^{\alpha} f\right|^{2} d x
$$

as an equivalent norm for $H_{n}\left(\mathbf{R}^{m}\right)$. With this interpretation in mind, it is not surprising that when we extend $D_{x}^{\alpha}$ to $H_{s}$, that $|\alpha| L^{2}$ derivatives are lost.

Lemma 1.1.3. $D_{x}^{\alpha}$ extends to define a continuous map $D_{x}^{\alpha}: H_{s} \rightarrow H_{s-|\alpha|}$.
Proof: Henceforth we will use $C$ to denote a generic constant. $C$ can depend upon certain auxiliary parameters which will usually be supressed in the interests of notational clarity. In this proof, for example, $C$ depends on $(s, \alpha)$ but not of course upon $f$. The estimate:

$$
\left|\xi^{\alpha}\right|^{2}\left(1+|\xi|^{2}\right)^{s-|\alpha|} \leq C\left(1+|\xi|^{2}\right)^{s}
$$

implies that:

$$
\left|D_{x}^{\alpha} f\right|_{s-\alpha}^{2}=\int\left|\xi^{\alpha} \hat{f}(\xi)\right|^{2}\left(1+|\xi|^{2}\right)^{s-|\alpha|} d \xi \leq C|f|_{s}^{2}
$$

for $f \in \mathcal{S}$. Since $H_{s}$ is the closure of $\mathcal{S}$ in the norm $\left.\right|_{s}$, this completes the proof.

We can also use the sup norm to measure derivatives. If $k$ is a nonnegative integer, we define:

$$
|f|_{\infty, k}=\sup _{x \in \mathbf{R}^{m}} \sum_{|\alpha| \leq k}\left|D_{x}^{\alpha} f\right| \quad \text { for } \quad f \in \mathcal{S}
$$

The completion of $\mathcal{S}$ with respect to this norm is a subset of $C^{k}\left(\mathbf{R}^{m}\right)$ (the continuous functions on $\mathbf{R}^{m}$ with continuous partial derivatives up to order $k)$. The next lemma relates the two norms $\left.\right|_{s}$ and $\left.\right|_{\infty, k}$. It will play an important role in showing the weak solutions we will construct to differential equations are in fact smooth.

Lemma 1.1.4. Let $k$ be a non-negative integer and let $s>k+\frac{m}{2}$. If $f \in H_{s}$, then $f$ is $C^{k}$ and there is an estimate $|f|_{\infty, k} \leq C|f|_{s}$. (Sobolev Lemma).

Proof: Suppose first $k=0$ and $f \in \mathcal{S}$. We compute

$$
\begin{aligned}
f(x) & =\int e^{i x \cdot \xi} \hat{f}(\xi) d \xi \\
& =\int\left\{e^{i x \cdot \xi} \hat{f}(\xi)\left(1+|\xi|^{2}\right)^{s / 2}\right\} \cdot\left\{\left(1+|\xi|^{2}\right)^{-s / 2}\right\} d \xi
\end{aligned}
$$

We apply the Cauchy-Schwarz inequality to estimate:

$$
|f(x)|^{2} \leq|f|_{s}^{2} \int\left(1+|\xi|^{2}\right)^{-s} d \xi
$$

Since $2 s>m,\left(1+|\xi|^{2}\right)^{-s}$ is integrable so $|f(x)| \leq C|f|_{s}$. We take the sup over $x \in \mathbf{R}^{m}$ to conclude $|f|_{\infty, 0} \leq C|f|_{s}$ for $f \in \mathcal{S}$. Elements of $H_{s}$ are the limits in the $\left.\right|_{s}$ norm of elements of $\mathcal{S}$. The uniform limit of continuous functions is continuous so the elements of $H_{s}$ are continuous and the same norm estimate extends to $H_{s}$. If $k>0$, we use the estimate:

$$
\left|D_{x}^{\alpha} f\right|_{\infty, 0} \leq C\left|D_{x}^{\alpha} f\right|_{s-|\alpha|} \leq C|f|_{s} \quad \text { for }|\alpha| \leq k \text { and } s-k>\frac{m}{2}
$$

to conclude $|f|_{\infty, k} \leq C|f|_{s}$ for $f \in \mathcal{S}$. A similar argument shows that the elements of $H_{s}$ must be $C^{k}$ and that this estimate continues to hold.

If $s>t$, we can estimate $\left(1+|\xi|^{2}\right)^{s} \geq\left(1+|\xi|^{2}\right)^{t}$. This implies that $|f|_{s} \geq|f|_{t}$ so the identity map on $\mathcal{S}$ extends to define an injection of $H_{s} \rightarrow H_{t}$ which is norm non-increasing. The next lemma shows that this injection is compact if we restrict the supports involved.

Lemma 1.1.5. Let $\left\{f_{n}\right\} \in \mathcal{S}$ be a sequence of functions with support in a fixed compact set $K$. We suppose there is a constant $C$ so $\left|f_{n}\right|_{s} \leq C$ for all $n$. Let $s>t$. There exists a subsequence $f_{n_{k}}$ which converges in $H_{t}$. (Rellich lemma).

Proof: Choose $g \in C_{0}\left(\mathbf{R}^{m}\right)$ which is identically 1 on a neighborhood of $K$. Then $g f_{n}=f_{n}$ so by Lemma 1.1.2(c) $\hat{f}_{n}=\hat{g} * \hat{f}_{n}$. We let $\partial_{j}=\frac{\partial}{\partial \xi_{j}}$ then $\partial_{j}\left(\hat{g} * \hat{f}_{n}\right)=\partial_{j} \hat{g} * \hat{f}_{n}$ so that:

$$
\left|\partial_{j} \hat{f}_{n}(\xi)\right| \leq \int\left|\left\{\partial_{j} \hat{g}(\xi-\zeta)\right\} \cdot \hat{f}_{n}(\zeta)\right| d \zeta
$$

We apply the Cauchy-Schwarz inequality to estimate:

$$
\left|\partial_{j} \hat{f}_{n}(\xi)\right| \leq\left|f_{n}\right|_{s} \cdot\left\{\int\left|\partial_{j} \hat{g}(\xi-\zeta)\right|^{2}\left(1+|\zeta|^{2}\right)^{-s} d \zeta\right\}^{1 / 2} \leq C \cdot h(\xi)
$$

where $h$ is some continuous function of $\xi$. A similar estimate holds for $\left|\hat{f}_{n}(\xi)\right|$. This implies that the $\left\{\hat{f}_{n}\right\}$ form a uniformly bounded equi-continuous family on compact $\xi$ subsets. We apply the Arzela-Ascoli theorem to extract a subsequence we again label by $f_{n}$ so that $\hat{f}_{n}(\xi)$ converges uniformly on compact subsets. We complete the proof by verifying that $f_{n}$ converges in $H_{t}$ for $s>t$. We compute:

$$
\left|f_{j}-f_{k}\right|_{t}^{2}=\int\left|\hat{f}_{j}-\hat{f}_{k}\right|^{2}\left(1+|\xi|^{2}\right)^{t} d \xi
$$

We decompose this integral into two parts, $|\xi| \geq r$ and $|\xi| \leq r$. On $|\xi| \geq r$ we estimate $\left(1+|\xi|^{2}\right)^{t} \leq\left(1+r^{2}\right)^{t-s}\left(1+|\xi|^{2}\right)^{s}$ so that:

$$
\begin{aligned}
\int_{|\xi| \geq r}\left|\hat{f}_{j}-\hat{f}_{k}\right|^{2}\left(1+|\xi|^{2}\right)^{t} d \xi & \leq\left(1+r^{2}\right)^{t-s} \int\left|\hat{f}_{j}-\hat{f}_{k}\right|^{2}\left(1+|\xi|^{2}\right)^{s} d \xi \\
& \leq 2 C\left(1+r^{2}\right)^{t-s}
\end{aligned}
$$

If $\varepsilon>0$ is given, we choose $r$ so that $2 C\left(1+r^{2}\right)^{t-s}<\varepsilon$. The remaining part of the integral is over $|\xi| \leq r$. The $\hat{f}_{j}$ converge uniformly on compact subsets so this integral can be bounded above by $\varepsilon$ if $j, k>j(\varepsilon)$. This completes the proof.

The hypothesis that the supports are uniformly bounded is essential. It is easy to construct a sequence $\left\{f_{n}\right\}$ with $\left|f_{n}\right|_{s}=1$ for all $n$ and such that the supports are pair-wise disjoint. In this case we can find $\varepsilon>0$ so that $\left|f_{j}-f_{k}\right|_{t}>\varepsilon$ for all $(j, k)$ so there is no convergent subsequence.

We fix $\phi \in \mathcal{S}$ and let $\phi_{\varepsilon}(x)=\phi(\varepsilon x)$. We suppose $\phi(0)=1$ and fix $f \in \mathcal{S}$. We compute:

$$
D_{x}^{\alpha}\left(f-\phi_{\varepsilon} f\right)=\left(1-\phi_{\varepsilon}\right) D_{x}^{\alpha} f+\text { terms of the form } \varepsilon^{j} D_{x}^{\beta} \phi(\varepsilon x) D_{x}^{\gamma} f
$$

As $\varepsilon \rightarrow 0$, these other terms go to zero in $L^{2}$. Since $\phi_{\varepsilon} \rightarrow 1$ pointwise, $\left(1-\phi_{\varepsilon}\right) D_{x} f$ goes to zero in $L^{2}$. This implies $\phi_{\varepsilon} f \rightarrow f$ in $H_{n}$ for any $n \geq 0$ as $\varepsilon \rightarrow 0$ and therefore $\phi_{\varepsilon} f \rightarrow f$ in $H_{s}$ for any $s$. If we take $\phi \in C_{0}^{\infty}\left(\mathbf{R}^{m}\right)$, this implies $C_{0}^{\infty}\left(\mathbf{R}^{m}\right)$ is dense in $H_{s}$ for any $s$.

Each $H_{s}$ space is a Hilbert space so it is isomorphic to its dual. Because there is no preferred norm for $H_{s}$, it is useful to obtain an invariant alternative characterization of the dual space $H_{s}^{*}$ :

Lemma 1.1.6. The $L^{2}$ pairing which maps $\mathcal{S} \times \mathcal{S} \rightarrow \mathbf{C}$ extends to a map of $H_{s} \times H_{-s} \rightarrow \mathbf{C}$ which is a perfect pairing and which identifies $H_{-s}$ with $H_{s}^{*}$. That is:
(a) $|(f, g)| \leq|f|_{s}|g|_{-s}$ for $f, g \in \mathcal{S}$,
(b) given $f \in \mathcal{S}$ there exists $g \in \mathcal{S}$ so $(f, g)=|f|_{s}|g|_{-s}$ and we can define

$$
|f|_{s}=\sup _{g \in \mathcal{S}, g \neq 0} \frac{|(f, g)|}{|g|_{-s}} .
$$

Proof: This follows from the fact that $H_{s}$ is $L^{2}$ with the weight function $\left(1+|\xi|^{2}\right)^{s}$ and $H_{-s}$ is $L^{2}$ with the weight function $\left(1+|\xi|^{2}\right)^{-s}$. We compute:

$$
(f, g)=(\hat{f}, \hat{g})=\int \hat{f}(\xi)\left(1+|\xi|^{2}\right)^{s / 2} \overline{\hat{g}}(\xi)\left(1+|\xi|^{2}\right)^{-s / 2} d \xi
$$

and apply the Cauchy-Schwartz inequality to prove (a).
To prove part (b), we note $|f|_{s} \geq \sup _{g \in \mathcal{S}, g \neq 0} \frac{|(f, g)|}{|g|_{-s}}$. We take $g$ to be defined by:

$$
\hat{g}=\hat{f}\left(1+|\xi|^{2}\right)^{s} \in \mathcal{S}
$$

and note that $(f, g)=(\hat{f}, \hat{g})=|f|_{s}^{2}$ and that $|g|_{-s}^{2}=|f|_{s}^{2}$ to see that equality can occur in (a) which proves (b)

If $s>t>u$ then we can estimate:

$$
(1+|\xi|)^{2 t} \leq \varepsilon(1+|\xi|)^{2 s}+C(\varepsilon)(1+|\xi|)^{2 u}
$$

for any $\varepsilon>0$. This leads immediately to the useful estimate:

Lemma 1.1.7. Let $s>t>u$ and let $\varepsilon>0$ be given. Then

$$
|f|_{t} \leq \varepsilon|f|_{s}+C(\varepsilon)|f|_{u}
$$

If $V$ is a finite dimensional vector space, let $C^{\infty}(V)$ be the space of smooth complex valued maps of $\mathbf{R}^{m} \rightarrow V$. We choose a fixed Hermitian inner product on $V$ and define $\mathcal{S}(V)$ and $H_{s}(V)$ as in the scalar case. If $\operatorname{dim}(V)=k$ and if we choose a fixed orthonormal basis for $V$, then $\mathcal{S}(V)$ and $H_{s}(V)$ become isomorphic to the direct sum of $k$ copies of $\mathcal{S}$ and of $H_{s}$. Lemmas 1.1.1 through 1.1.7 extend in the obvious fashion.

We conclude this subsection with an extremely useful if elementary estimate:

Lemma 1.1.8. (Peetre's Inequality). Let $s$ be real and $x, y \in \mathbf{R}^{m}$. Then $(1+|x+y|)^{s} \leq(1+|y|)^{s}(1+|x|)^{|s|}$.
Proof: We suppose first $s>0$. We raise the triangle inequality:

$$
1+|x+y|<1+|x|+|y| \leq(1+|y|)(1+|x|)
$$

to the $s^{\text {th }}$ power to deduce the desired inequality. We now suppose $s<0$. A similar inequality:

$$
(1+|y|)^{-s} \leq(1+|x+y|)^{-s}(1+|x|)^{-s}
$$

yields immediately:

$$
(1+|x+y|)^{s} \leq(1+|y|)^{s}(1+|x|)^{-s}
$$

to complete the proof.

### 1.2. Pseudo-Differential Operators on $\mathbf{R}^{m}$.

A linear partial differential operator of order $d$ is a polynomial expression $P=p(x, D)=\sum_{|\alpha| \leq d} a_{\alpha}(x) D_{x}^{\alpha}$ where the $a_{\alpha}(x)$ are smooth. The symbol $\sigma P=p$ is defined by:

$$
\sigma P=p(x, \xi)=\sum_{|\alpha| \leq d} a_{\alpha}(x) \xi^{\alpha}
$$

and is a polynomial of order $d$ in the dual variable $\xi$. It is convenient to regard the pair $(x, \xi)$ as defining a point of the cotangent space $T^{*}\left(\mathbf{R}^{m}\right)$; we will return to this point again when we discuss the effect of coordinate transformations. The leading symbol $\sigma_{L} P$ is the highest order part:

$$
\sigma_{L} P(x, \xi)=\sum_{|\alpha|=d} a_{\alpha}(x) \xi^{\alpha}
$$

and is a homogeneous polynomial of order $d$ in $\xi$.
We can use the Fourier inversion formula to express:

$$
P f(x)=\int e^{i x \cdot \xi} p(x, \xi) \hat{f}(\xi) d \xi=\int e^{i(x-y) \cdot \xi} p(x, \xi) f(y) d y d \xi
$$

for $f \in \mathcal{S}$. We note that since the second integral does not converge absolutely, we cannot interchange the $d y$ and $d \xi$ orders of integration. We use this formalism to define the action of pseudo-differential operators ( $\Psi$ DO's) for a wider class of symbols $p(x, \xi)$ than polynomials. We make the following
Definition. $p(x, \xi)$ is a symbol of order $d$ and we write $p \in S^{d}$ if
(a) $p(x, \xi)$ is smooth in $(x, \xi) \in \mathbf{R}^{m} \times \mathbf{R}^{m}$ with compact $x$ support,
(b) for all $(\alpha, \beta)$ there are constants $C_{\alpha, \beta}$ such that

$$
\left|D_{x}^{\alpha} D_{\xi}^{\beta} p(x, \xi)\right| \leq C_{\alpha, \beta}(1+|\xi|)^{d-|\beta|}
$$

For such a symbol $p$, we define the associated operator $P(x, D)$ by:

$$
P(x, D)(f)(x)=\int e^{i x \cdot \xi} p(x, \xi) \hat{f}(\xi) d \xi=\int e^{(x-y) \cdot \xi} p(x, \xi) f(y) d y d \xi
$$

as a linear operator mapping $\mathcal{S} \rightarrow \mathcal{S}$.
A differential operator has as its order a positive integer. The order of a pseudo-differential operator is not necessarily an integer. For example, if $f \in C_{0}^{\infty}\left(\mathbf{R}^{m}\right)$, define:

$$
p(x, \xi)=f(x)\left(1+|\xi|^{2}\right)^{d / 2} \in S^{d} \quad \text { for any } d \in \mathbf{R} .
$$

This will be a symbol of order $d$. If $p \in S^{d}$ for all $d$, then we say that $p \in S^{-\infty}$ is infinitely smoothing. We adopt the notational convention of letting $p, q, r$ denote symbols and $P, Q, R$ denote the corresponding $\Psi$ DO's.

Because we shall be interested in problems on compact manifolds, we have assumed the symbols have compact $x$ support to avoid a number of technical complications. The reader should note that there is a well defined theory which does not require compact $x$ support.

When we discuss the heat equation, we shall have to consider a wider class of symbols which depend on a complex parameter. We postpone discussion of this class until later to avoid unnecessarily complicating the discussion at this stage. We shall phrase the theorems and proofs of this section in such a manner that they will generalize easily to the wider class of symbols.

Our first task is to extend the action of $P$ from $\mathcal{S}$ to $H_{s}$.
Lemma 1.2.1. Let $p \in S^{d}$ then $|P f|_{s-d} \leq C|f|_{s}$ for $f \in \mathcal{S}$. $P$ extends to a continuous map $P: H_{s} \rightarrow H_{s-d}$ for all $s$.
Proof: We compute $P f(x)=\int e^{i x \cdot \xi} p(x, \xi) \hat{f}(\xi) d \xi$ so that the Fourier transform is given by:

$$
\widehat{P f}(\zeta)=\int e^{i x \cdot(\xi-\zeta)} p(x, \xi) \hat{f}(\xi) d \xi d x
$$

This integral is absolutely convergent since $p$ has compact $x$ support so we may interchange the order of integration. If we define

$$
q(\zeta, \xi)=\int e^{-i x \cdot \zeta} p(x, \xi) d x
$$

as the Fourier transform in the $x$ direction, then

$$
\widehat{P f}(\zeta)=\int q(\zeta-\xi, \xi) \hat{f}(\xi) d \xi
$$

By Lemma 1.1.6, $|P f|_{s-d}=\sup _{g \in S} \frac{|(P f, g)|}{|g|_{d-s}}$. We compute:

$$
(P f, g)=\int q(\zeta-\xi, \xi) \hat{f}(\xi) \overline{\hat{g}}(\zeta) d \xi d \zeta
$$

Define:

$$
K(\zeta, \xi)=q(\zeta-\xi, \xi)(1+|\xi|)^{-s}(1+|\zeta|)^{s-d}
$$

then:

$$
(P f, g)=\int K(\zeta, \xi) \hat{f}(\xi)(1+|\xi|)^{s} \overline{\hat{g}}(\zeta)(1+|\zeta|)^{d-s} d \xi d \zeta
$$

We apply the Cauchy-Schwarz inequality to estimate:

$$
\begin{aligned}
&|(P f, g)| \leq\left\{\int|K(\zeta, \xi)||\hat{f}(\xi)|^{2}(1+|\xi|)^{2 s} d \xi d \zeta\right\}^{1 / 2} \\
& \times\left\{\int|K(\zeta, \xi)||\hat{g}(\zeta)|^{2}(1+|\zeta|)^{2 d-2 s} d \xi d \zeta\right\}^{1 / 2}
\end{aligned}
$$

We complete the proof by showing

$$
\int|K(\zeta, \xi)| d \xi \leq C \quad \text { and } \quad \int|K(\zeta, \xi)| d \zeta \leq C
$$

since then $|(P f, g)| \leq C|f|_{s}|g|_{d-s}$.
By hypothesis, $p$ has suppport in a compact set $\mathcal{K}$ and we have estimates:

$$
\left|D_{x}^{\alpha} p(x, \xi)\right| \leq C_{\alpha}(1+|\xi|)^{d}
$$

Therefore:

$$
\left|\zeta^{\alpha} q(\zeta, \xi)\right|=\left|\int e^{-i x \cdot \zeta} D_{x}^{\alpha} p(x, \xi) d x\right| \leq C_{\alpha}(1+|\xi|)^{d} \operatorname{vol}(\mathcal{K})
$$

Therefore, for any integer $k,|q(\zeta, \xi)| \leq C_{k}(1+|\xi|)^{d}(1+|\zeta|)^{-k} \operatorname{vol}(\mathcal{K})$ and:

$$
|K(\zeta, \xi)| \leq C_{k}(1+|\xi|)^{d-s}(1+|\zeta|)^{s-d}(1+|\zeta-\xi|)^{-k} \operatorname{vol}(\mathcal{K})
$$

We apply Lemma 1.1.8 with $x+y=\xi$ and $y=\zeta$ to estimate:

$$
|K(\zeta, \xi)| \leq C_{k}(1+|\zeta-\xi|)^{|d-s|-k} \operatorname{vol}(\mathcal{K})
$$

If we choose $k>\frac{m}{2}+|d-s|$, then this will be integrable and complete the proof.

Our next task is to show that the class of $\Psi$ DO's forms an algebra under the operations of composition and taking adjoint. Before doing that, we study the situation with respect to differential operators to motivate the formulas we shall derive. Let $P=\sum_{\alpha} p_{\alpha}(x) D_{x}^{\alpha}$ and let $Q=\sum_{\alpha} q_{\alpha}(x) D_{x}^{\alpha}$ be two differential operators. We assume $p$ and $q$ have compact $x$ support. It is immediate that:

$$
P^{*}=\sum_{\alpha} D_{x}^{\alpha} p_{\alpha}^{*} \quad \text { and } \quad P Q=\sum_{\alpha, \beta} p_{\alpha}(x) D_{x}^{\alpha} q_{\beta}(x) D_{x}^{\beta}
$$

are again differential operators in our class. Furthermore, using Leibnitz's rule

$$
\begin{aligned}
D_{x}^{\alpha}(f g) & =\sum_{\beta+\gamma=\alpha} D_{x}^{\beta}(f) \cdot D_{x}^{\gamma}(g) \cdot \frac{\alpha!}{\beta!\gamma!}, \\
d_{\xi}^{\beta}\left(\xi^{\beta+\gamma}\right) & =\xi^{\gamma} \cdot \frac{(\beta+\gamma)!}{\gamma!},
\end{aligned}
$$

it is an easy combinatorial exercise to compute that:

$$
\sigma\left(P^{*}\right)=\sum_{\alpha} d_{\xi}^{\alpha} D_{x}^{\alpha} p^{*} / \alpha!\quad \text { and } \quad \sigma(P Q)=\sum_{\alpha} d_{\xi}^{\alpha} p \cdot D_{x}^{\alpha} q / \alpha!
$$

The perhaps surprising fact is that these formulas remain true in some sense for $\Psi D O$ 's, only the sums will become infinite rather than finite.

We introduce an equivalence relation on the class of symbols by defining $p \sim q$ if $p-q \in S^{-\infty}$. We note that if $p \in S^{-\infty}$ then $P: H_{s} \rightarrow H_{t}$ for all $s$ and $t$ by Lemma 1.2.1. Consequently by Lemma 1.1.4, $P: H_{s} \rightarrow C_{0}^{\infty}$ for all $s$ so that $P$ is infinitely smoothing in this case. Thus we mod out by infinitely smoothing operators.

Given symbols $p_{j} \in S^{d_{j}}$ where $d_{j} \rightarrow-\infty$, we write

$$
p \sim \sum_{j=1}^{\infty} p_{j}
$$

if for every $d$ there is an integer $k(d)$ such that $k \geq k(d)$ implies that $p-\sum_{j=1}^{k} p_{j} \in S^{d}$. We emphasize that this sum does not in fact need to converge. The relation $p \sim \sum p_{j}$ simply means that the difference between $P$ and the partial sums of the $P_{j}$ is as smoothing as we like. It will turn out that this is the appropriate sense in which we will generalize the formulas for $\sigma\left(P^{*}\right)$ and $\sigma(P Q)$ from differential to pseudo-differential operators.

Ultimately, we will be interested in operators which are defined on compact manifolds. Consequently, it poses no difficulties to restrict the domain and the range of our operators. Let $U$ be a open subset of $\mathbf{R}^{m}$ with compact closure. Let $p(x, \xi) \in S^{d}$ have $x$ support in $U$. We restrict the domain of the operator $P$ to $C_{0}^{\infty}(U)$ so $P: C_{0}^{\infty}(U) \rightarrow C_{0}^{\infty}(U)$. Let $\Psi_{d}(U)$ denote the space of all such operators. For $d \leq d^{\prime}$, then $\Psi_{d}(U) \subseteq \Psi_{d^{\prime}}(U)$. We define

$$
\Psi(U)=\bigcup_{d} \Psi_{d}(U) \quad \text { and } \quad \Psi_{-\infty}(U)=\bigcap_{d} \Psi_{d}(U)
$$

to be the set of all pseudo-differential operators on $U$ and the set of infinitely smoothing pseudo-differential operators on $U$.

More generally, let $p(x, \xi)$ be a matrix valued symbol; we suppose the components of $p$ all belong to $S^{d}$. The corresponding operator $P$ is given by a matrix of pseudo-differential operators. $P$ is a map from vector valued functions with compact support in $U$ to vector valued functions with compact support in $U$. We shall not introduce separate notation for the shape of $p$ and shall continue to denote the collection of all such operators by $\Psi_{d}(U)$. If $p$ and $q$ are matrix valued and of the proper shape, we define $p \cdot q$ and also the operator $P \cdot Q$ by matrix product and by composition. We also define $p^{*}$ and $P^{*}$ to be the matrix adjoint and the operator adjoint so that $\left(P^{*} f, g\right)=\left(f, P^{*} g\right)$ where $f$ and $g$ are vector valued and of compact support. Before studying the algebra structure on $\Psi(U)$, we must enlarge the class of symbols which we can admit:

Lemma 1.2.2. Let $r(x, \xi, y)$ be a matrix valued symbol which is smooth in $(x, \xi, y)$. We suppose $r$ has compact $x$ support inside $U$ and that there are estimates:

$$
\left|D_{x}^{\alpha} D_{\xi}^{\beta} D_{y}^{\gamma} r\right| \leq C_{\alpha, \beta, \gamma}(1+|\xi|)^{d-|\beta|}
$$

for all multi-indices $(\alpha, \beta, \gamma)$. If $f$ is vector valued with compact support in $U$, we define:

$$
R f(x)=\int e^{i(x-y) \cdot \xi} r(x, \xi, y) f(y) d y d \xi
$$

Then this operator is in $\Psi_{d}(U)$ and the symbol is given by:

$$
\left.\sigma R(x, \xi) \sim\left\{\sum_{\alpha} d_{\xi}^{\alpha} D_{y}^{\alpha} r(x, \xi, y) / \alpha!\right\}\right|_{x=y}
$$

Proof: We note that any symbol in $S^{d}$ belongs to this class of operators if we define $r(x, \xi, y)=p(x, \xi)$. We restricted to vector valued functions with compact support in $U$. By multiplying $r$ by a cut-off function in $y$ with compact support which is 1 over $U$, we may assume without loss of generality the $y$ support of $r$ is compact as well. Define:

$$
q(x, \xi, \zeta)=\int e^{-i y \cdot \zeta} r(x, \xi, y) d y
$$

to be the Fourier transform of $r$ in the $y$ variable. Using Lemma 1.1.2 we see $\widehat{(r f)}=\hat{r} * \hat{f}$. This implies that:

$$
\int e^{-i y \cdot \xi} r(x, \xi, y) f(y) d y=\int q(x, \xi, \xi-\zeta) \hat{f}(\zeta) d \zeta
$$

The argument given in the proof of Lemma 1.2.1 gives estimates of the form:

$$
|q(x, \xi, \zeta)| \leq C_{k}(1+|\xi|)^{d}(1+|\zeta|)^{-k} \quad \text { and } \quad|\hat{f}(\zeta)| \leq C_{k}(1+|\zeta|)^{-k}
$$

for any $k$. Consequently:

$$
|q(x, \xi, \xi-\zeta) \hat{f}(\zeta)| \leq C_{k}(1+|\xi|)^{d}(1+|\xi-\zeta|)^{-k}(1+|\zeta|)^{-k}
$$

We apply Lemma 1.1.8 to estimate:

$$
|q(x, \xi, \xi-\zeta) \hat{f}(\zeta)| \leq C_{k}(1+|\xi|)^{|d|-k}(1+|\zeta|)^{|d|-k}
$$

so this is absolutely integrable. We change the order of integration and express:

$$
R f(x)=\int e^{i x \cdot \xi} q(x, \xi, \xi-\zeta) \hat{f}(\zeta) d \xi d \zeta
$$

We define:

$$
p(x, \zeta)=\int e^{i x(\xi-\zeta)} q(x, \xi, \xi-\zeta) d \xi
$$

and compute:

$$
R f(x)=\int e^{i x \cdot \zeta} p(x, \zeta) \hat{f}(\zeta) d \zeta
$$

is a pseudo-differential operator once it is verified that $p(x, \zeta)$ is a symbol in the correct form.

We change variables to express:

$$
p(x, \zeta)=\int e^{i x \cdot \xi} q(x, \xi+\zeta, \xi) d \xi
$$

and estimate:

$$
\begin{aligned}
|q(x, \xi+\zeta, \xi)| & \leq C_{k}(1+|\xi+\zeta|)^{d}(1+|\xi|)^{-k} \\
& \leq C_{k}(1+|\zeta|)^{d}(1+|\xi|)^{|d|-k}
\end{aligned}
$$

This is integrable so $|p(x, \zeta)| \leq C_{k}^{\prime}(1+|\zeta|)^{d}$. Similar estimates on $\mid D_{x}^{\alpha} D_{\zeta}^{\beta}$ $q(x, \xi+\zeta, \xi) \mid$ which arise from the given estimates for $r$ show that $p \in S^{d}$ so that $R$ is a pseudo-differential operator.

We use Taylor's theorem on the middle variable of $q(x, \xi+\zeta, \xi)$ to expand:

$$
q(x, \xi+\zeta, \xi)=\sum_{|\alpha| \leq k} \frac{d_{\zeta}^{\alpha} q(x, \zeta, \xi) \xi^{\alpha}}{\alpha!}+q_{k}(x, \zeta, \xi)
$$

The remainder $q_{k}$ decays to arbitrarily high order in $(\xi, \zeta)$ and after integration gives rise to a symbol in $S^{d-k}$ which may therefore be ignored. We integrate to conclude

$$
\begin{aligned}
p(x, \zeta) & =\sum_{|\alpha| \leq k} \int e^{i x \cdot \xi} \frac{d_{\zeta}^{\alpha} q(x, \zeta, \xi) \xi^{\alpha}}{\alpha!} d \xi+\text { remainder } \\
& =\left.\sum_{|\alpha| \leq k} \frac{d_{\zeta}^{\alpha} D_{y}^{\alpha} r(x, \zeta, y)}{\alpha!}\right|_{x=y}+\text { a remainder }
\end{aligned}
$$

using Lemma 1.1.2. This completes the proof of the lemma.
We use this technical lemma to show that the pseudo-differential operators form an algebra:

Lemma 1.2.3. Let $P \in \Psi_{d}(U)$ and let $Q \in \Psi_{e}(U)$. Then:
(a) If $U^{\prime}$ is any open set with compact closure containing $\bar{U}$, then $P^{*} \in$ $\Psi_{d}\left(U^{\prime}\right)$ and $\sigma\left(P^{*}\right) \sim \sum_{\alpha} d_{\xi}^{\alpha} D_{x}^{\alpha} p^{*} / \alpha!$.
(b) Assume that $P$ and $Q$ have the proper shapes so $P Q$ and $p q$ are defined. Then $P Q \in \Psi_{d+e}(U)$ and $\sigma(P Q) \sim \sum_{\alpha} d_{\xi}^{\alpha} p \cdot D_{x}^{\alpha} q / \alpha$ !.

Proof: The fact that $P^{*}$ lies in a larger space is only a slight bit of technical bother; this fact plays an important role in considering boundary value problems of course. Let $(f, g)=f \cdot g$ be the pointwise Hermitian inner product. Fix $\phi \in C_{0}^{\infty}\left(U^{\prime}\right)$ to be identically 1 on $U$ and compute:

$$
\begin{aligned}
(P f, g) & =\int e^{i(x-y) \cdot \xi} p(x, \xi) \phi(y) f(y) \cdot g(x) d y d \xi d x \\
& =\int f(y) \cdot e^{i(y-x) \cdot \xi} p^{*}(x, \xi) \phi(y) g(x) d y d \xi d x
\end{aligned}
$$

since the inner product is Hermitian. By approximating $p^{*}(x, \xi)$ by functions with compact $\xi$ support, we can justify the use of Fubini's theorem to replace $d y d \xi d x$ by $d x d \xi d y$ and express:

$$
\begin{aligned}
(P f, g) & =\int f(y) \cdot e^{i(y-x) \cdot \xi} p^{*}(x, \xi) \phi(y) g(x) d x d \xi d y \\
& =\left(f, P^{*} g\right)
\end{aligned}
$$

where we define:

$$
P^{*} g(y)=\int e^{i(y-x) \cdot \xi} p^{*}(x, \xi) \phi(y) g(x) d x d \xi
$$

This is an operator of the form discussed in Lemma 1.2.2 so $P^{*} \in \Psi_{d}\left(U^{\prime}\right)$ and we compute:

$$
\sigma\left(P^{*}\right) \sim \sum_{\alpha} d_{\xi}^{\alpha} D_{x}^{\alpha} p^{*} / \alpha!
$$

since $\phi=1$ on the support of $p$. This completes the proof of (a). We note that we can delete the factor of $\phi$ from the expression for $P^{*} g$ since it was only needed to prove $P^{*}$ was a $\Psi D O$.

We use (a) to prove (b). Since:

$$
Q^{*} g(y)=\int e^{i(y-x) \cdot \xi} q^{*}(x, \xi) g(x) d x d \xi
$$

the Fourier inversion formula implies:

$$
\left(\widehat{Q^{*} g}\right)=\int e^{-i x \cdot \xi} q^{*}(x, \xi) g(x) d x
$$

If $\tilde{q}$ is the symbol of $Q^{*}$, then we interchange the roles of $Q$ and $Q^{*}$ to see:

$$
\widehat{(Q g})=\int e^{-i x \cdot \xi} \tilde{q}^{*}(x, \xi) g(x) d x
$$

Therefore:

$$
\begin{aligned}
P Q g(x) & \left.=\int e^{i x \cdot \xi} p(x, \xi) \widehat{(Q g}\right)(\xi) d \xi \\
& =\int e^{i(x-y) \cdot \xi} p(x, \xi) \tilde{q}^{*}(y, \xi) g(y) d y d \xi
\end{aligned}
$$

which is an operator of the form discussed in Lemma 1.2.2 if $r(x, \xi, y)=$ $p(x, \xi) \tilde{q}^{*}(y, \xi)$. This proves $P Q$ is a pseudo-differential operator of the correct order. We compute the symbol of $P Q$ to be:

$$
\sim \sum_{\alpha} d_{\xi}^{\alpha} D_{y}^{\alpha}\left(p(x, \xi) \tilde{q}^{*}(y, \xi)\right) / \alpha!\quad \text { evaluated at } x=y
$$

We use Leibnitz's formula and expand this in the form:

$$
\sim \sum_{\beta, \gamma} d_{\xi}^{\beta} p(x, \xi) D_{y}^{\beta} d_{\xi}^{\gamma} D_{y}^{\gamma} \tilde{q}^{*} / \beta!\gamma!
$$

The sum over $\gamma$ yields the symbol of $Q^{* *}=Q$ so we conclude finally

$$
\sigma(P Q) \sim \sum_{\beta} d_{\xi}^{\beta} p(x, \xi) D_{x}^{\beta} q(x, \xi) / \beta!
$$

which completes the proof.
Let $K(x, y)$ be a smooth matrix valued function with compact $x$ support in $U$. If $f$ is vector valued with compact support in $U$, we define:

$$
P(K)(f)(x)=\int K(x, y) f(y) d y
$$

Lemma 1.2.4. Let $K(x, y)$ be smooth with compact $x$ support in $U$, then $P(K) \in \Psi_{-\infty}(U)$.

Proof: We let $\phi(\xi) \in C_{0}^{\infty}\left(\mathbf{R}^{m}\right)$ with $\int \phi(\xi) d \xi=1$. Define:

$$
r(x, \xi, y)=e^{i(y-x) \cdot \xi} \phi(\xi) K(x, y)
$$

then this is a symbol in $S^{-\infty}$ of the sort discussed in Lemma 1.2.2. It defines an infinitely smoothing operator. It is immediate that:

$$
P(K)(f)(x)=\int e^{i(x-y) \cdot \xi} r(x, \xi, y) f(y) d y d \xi
$$

Conversely, it can be shown that any infinitely smoothing map has a smooth kernel. In general, of course, it is not possible to represent an arbitrary pseudo-differential operator by a kernel. If $P$ is smoothing enough, however, we can prove:

Lemma 1.2.5. Let $r$ satisfy the hypothesis of Lemma 1.2 .2 where $d<$ $-m-k$. We define $K(x, y)=\int e^{i(x-y) \cdot \xi} r(x, \xi, y) d \xi$. Then $K$ is $C^{k}$ in $(x, y)$ and $R f(x)=\int K(x, y) f(y) d y$.

Proof: If we can show $K$ is well defined, then the representation of $R$ in terms of the kernel $K$ will follow from Fubini's theorem. We estimate:

$$
\begin{aligned}
D_{x}^{\alpha} D_{y}^{\beta} K(x, y)= & \sum_{\substack{\alpha=\alpha_{1}+\alpha_{2} \\
\beta=\beta_{1}+\beta_{2}}} \frac{\alpha!\beta!}{\alpha_{1}!\alpha_{2}!\beta_{1}!\beta_{2}!}(-1)^{\left|\beta_{1}\right|} \\
& \times\left\{\int e^{i(x-y) \cdot \xi} \xi^{\alpha_{1}+\beta_{1}} D_{x}^{\alpha_{2}} D_{y}^{\beta_{2}} r(x, \xi, y) d \xi\right\} .
\end{aligned}
$$

Since we can estimate:

$$
\left|\xi^{\alpha_{1}+\beta_{1}} D_{x}^{\alpha_{2}} D_{y}^{\beta_{2}} r(x, \xi, y)\right| \leq C(1+|\xi|)^{d+|\alpha|+|\beta|}
$$

this will be integrable for $|\alpha|+|\beta| \leq k$. Thus $K$ is $C^{k}$ and the representation of $R$ follows immediately.

In Lemma 1.2.2 we computed the symbol of the pseudo-differential operator defined by $r(x, \xi, y)$ in terms of $d_{\xi}^{\alpha} D_{y}^{\alpha} r$ when $x=y$. This implies the singular (i.e., the non-smoothing part) of $R$ is concentrated near the diagonal $x=y$. We make this more precise:

Lemma 1.2.6. Let $r(x, \xi, y)$ satisfy the hypothesis of Lemma 1.2.2. Suppose the $x$ support of $r$ is disjoint from the $y$ support of $r$, then $R$ is infinitely smoothing and is represented by a smooth kernel function $K(x, y)$.

Proof: We would like to define $K(x, y)=\int e^{(x-y) \cdot \xi} r(x, \xi, y) d \xi$. Unfortunately, this integral need not converge in general. By hypothesis, $|x-y| \geq \varepsilon>0$ on the support of $r$. We define the Laplacian $\Delta_{\xi}=\sum_{\nu} D_{\xi_{\nu}}^{2}$. Since $\Delta_{\xi} e^{i(x-y) \cdot \xi}=|x-y|^{2} e^{i(x-y) \cdot \xi}$ we integrate by parts in a formal sense $k$ times to express:

$$
R f(x)=\int e^{i(x-y) \cdot \xi}|x-y|^{-2 k} \Delta_{\xi}^{k} r(x, \xi, y) f(y) d y d \xi
$$

This formal process may be justified by first approximating $r$ by a function with compact $\xi$ support. We now define

$$
K(x, y)=\int e^{i(x-y) \cdot \xi}|x-y|^{-2 k} \Delta_{\xi}^{k} r(x, \xi, y) d \xi
$$

for any $k$ sufficiently large. Since $\Delta_{\xi}^{k} r$ decays to arbitrarily high order in $\xi$, we use the same argument as that given in Lemma 1.2.5 to show that
$K(x, y)$ is arbitrarily smooth in $(x, y)$ and hence is $C^{\infty}$. This completes the proof.

We note that in general $K(x, y)$ will become singular at $x=y$ owing to the presence of the terms $|x-y|^{-2 k}$ if we do not assume the support of $x$ is disjoint from the support of $y$.

A differential operator $P$ is local in the sense that if $f=0$ on some open subset of $U$, then $P f=0$ on that same subset since differentiation is a purely local process. $\Psi$ DO's are not local in general since they are defined by the Fourier transform which smears out the support. Nevertheless, they do have a somewhat weaker property, they do not smear out the singular support of a distribution $f$. More precisely, let $f \in H_{s}$. If $\phi \in C_{0}^{\infty}(U)$, we define the map $f \mapsto \phi f$. If we take $r(x, \xi, y)=\phi(x)$ and apply Lemma 1.2.2, then we see that this is a pseudo-differential operator of order 0. Therefore $\phi f \in H_{s}$ as well. This gives a suitable notion of restriction. We say that $f$ is smooth on an open subset $U^{\prime}$ of $U$ if and only if $\phi f \in C^{\infty}$ for every such $\phi$. An operator $P$ is said to be pseudo-local if $f$ is smooth on $U^{\prime}$ implies $P f$ is smooth on $U^{\prime}$.

Lemma 1.2.7. Pseudo-differential operators are pseudo-local.
Proof: Let $P \in \Psi_{d}(U)$ and let $f \in H_{s}$. Fix $x \in U^{\prime}$ and choose $\phi \in$ $C_{0}^{\infty}\left(U^{\prime}\right)$ to be identically 1 near $x$. Choose $\psi \in C_{0}^{\infty}\left(U^{\prime}\right)$ with support contained in the set where $\phi$ is identically 1 . We must verify that $\psi P f$ is smooth. We compute:

$$
\psi P f=\psi P \phi f+\psi P(1-\phi) f
$$

By hypothesis, $\phi f$ is smooth so $\psi P \phi f$ is smooth. The operator $\psi P(1-\phi)$ is represented by a kernel of the form $\psi(x) p(x, \xi)(1-\phi(y))$ which has disjoint $x$ and $y$ support. Lemma 1.2.6 implies $\psi P(1-\phi) f$ is smooth which completes the proof.

In Lemmas 1.2.2 and 1.2.3 we expressed the symbol of an operator as a infinite asymptotic series. We show that the algebra of symbols is complete in a certain sense:

Lemma 1.2.8. Let $p_{j} \in S^{d_{j}}(U)$ where $d_{j} \rightarrow-\infty$. Then there exists $p \sim \sum_{j} p_{j}$ which is a symbol in our class. $p$ is a unique modulo $S^{-\infty}$.
Proof: We may assume without loss of generality that $d_{1}>d_{2}>\cdots \rightarrow$ $-\infty$. We will construct $p \in S^{d_{1}}$. The uniqueness is clear so we must prove existence. The $p_{j}$ all have support inside $U$; we will construct $p$ with support inside $U^{\prime}$ where $U^{\prime}$ is any open set containing the closure of $U$.

Fix a smooth function $\phi$ such that:

$$
0 \leq \phi \leq 1, \quad \phi(\xi)=0 \text { for }|\xi| \leq 1, \quad \phi(\xi)=1 \text { for }|\xi| \geq 2
$$

We use $\phi$ to cut away the support near $\xi=0$. Let $t_{j} \rightarrow 0$ and define:

$$
p(x, \xi)=\sum_{j} \phi\left(t_{j} \xi\right) p_{j}(x, \xi) .
$$

For any fixed $\xi, \phi\left(t_{j} \xi\right)=0$ for all but a finite number of $j$ so this sum is well defined and smooth in $(x, \xi)$. For $j>1$ we have

$$
\left|p_{j}(x, \xi)\right| \leq C_{j}(1+|\xi|)^{d_{j}}=C_{j}(1+|\xi|)^{d_{1}}(1+|\xi|)^{d_{j}-d_{1}} .
$$

If $|\xi|$ is large enough, $(1+|\xi|)^{d_{j}-d_{1}}$ is as small as we like and therefore by passing to a subsequence of the $t_{j}$ we can assume

$$
\left|\phi\left(t_{j} \xi\right) p_{j}(x, \xi)\right| \leq 2^{-j}(1-|\xi|)^{d_{1}} \quad \text { for } j>1
$$

This implies that $|p(x, \xi)| \leq\left(C_{1}+1\right)(1+|\xi|)^{d_{1}}$. We use a similar argument with the derivatives and use a diagonalization argument on the resulting subsequences to conclude $p \in S^{d}$. The supports of all the $p_{j}$ are contained compactly in $U$ so the support of $p$ is contained in $\bar{U}$ which is contained in $U^{\prime}$.

We now apply exactly the same argument to $p_{d_{2}}+\cdots$ to assume that $p_{d_{2}}+\cdots \in S^{d_{2}}$. We continue in this fashion and use a diagonalization argument on the resulting subsequences to conclude in the end that

$$
\sum_{j=j_{0}}^{\infty} \phi\left(t_{j} \xi\right) p_{j}(x, \xi) \in S^{k} \quad \text { for } k=d_{j_{0}}
$$

Since $p_{j}-\phi\left(t_{j} \xi\right) p_{j} \in S^{-\infty}$, this implies $p-\sum_{j=1}^{j_{0}} p_{j} \in S^{k}$ and completes the proof.

If $K(x, y)$ is smooth with compact $x, y$ support in $U$, then $P(K) \in$ $\Psi_{-\infty}(U)$ defines a continuous operator from $H_{s} \rightarrow H_{t}$ for any $s, t$. Let $|P|_{s, t}$ denote the operator norm so $|P f|_{t} \leq|P|_{s, t}|f|_{s}$ for any $f \in \mathcal{S}$. It will be convenient to be able to estimate $|K|_{\infty, k}$ in terms of these norms:

Lemma 1.2.9. Let $K(x, y)$ be a smooth kernel with compact $x, y$ support in $U$. Let $P=P(K)$ be the operator defined by $K$. If $k$ is a non-negative integer, then $|K|_{\infty, k} \leq C(k)|P|_{-k, k}$

Proof: By arguing separately on each entry in the matrix $K$, we may reduce ourselves to the scalar case. Suppose first $k=0$. Choose $\phi \in$ $C_{0}^{\infty}\left(\mathbf{R}^{m}\right)$ positive with $\int \phi(x) d x=1$. Fix points $\left(x_{0}, y_{0}\right) \in U \times U$ and define:

$$
f_{n}(x)=n^{m} \phi\left(n\left(x-x_{0}\right)\right) \quad \text { and } \quad g_{n}(y)=n^{m} \phi\left(n\left(y-y_{0}\right)\right) .
$$

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Then if $n$ is large, $f_{n}$ and $g_{n}$ have compact support in $U$. Then:

$$
K\left(x_{0}, y_{0}\right)=\lim _{n \rightarrow \infty} \int f_{n}(x) K(x, y) g_{n}(y) d y d x=\lim _{n \rightarrow \infty}\left(f_{n}, P g_{n}\right)
$$

by Lemma 1.1.1. We estimate

$$
\left|\left(f_{n}, P g_{n}\right)\right| \leq|P|_{0,0}\left|f_{n}\right|_{0}\left|g_{n}\right|_{0}=|P|_{0,0}|\phi|_{0}^{2}
$$

to complete the proof in this case.
If $|\alpha| \leq k,|\beta| \leq k$ then:

$$
\begin{aligned}
D_{x}^{\alpha} D_{y}^{\beta} K(x, y) & =\lim _{n \rightarrow \infty} \int f_{n}(x)\left\{D_{x}^{\alpha} D_{y}^{\beta} K(x, y)\right\} g_{n}(y) d y d x \\
& =\lim _{n \rightarrow \infty} \int\left(D_{x}^{\alpha} f_{n}\right) K(x, y)\left(D_{y}^{\beta} g_{n}(y)\right) d y d x \\
& =\lim _{n \rightarrow \infty}\left(D_{x}^{\alpha} f_{n}, P D_{y}^{\beta} g_{n}\right) .
\end{aligned}
$$

We use Lemma 1.1.6 to estimate this by

$$
\begin{aligned}
\left|D_{x}^{\alpha} f_{n}\right|_{-k}\left|P D_{y}^{\beta} g_{n}\right|_{k} & \leq\left|f_{n}\right|_{0}|P|_{-k, k}\left|D_{y}^{\beta} g_{n}\right|_{-k} \\
& \leq\left|f_{n}\right|_{0}\left|g_{n}\right|_{0}|P|_{-k, k}=|P|_{-k, k}|\phi|_{0}^{2}
\end{aligned}
$$

to complete the proof.

### 1.3. Ellipticity and Pseudo-Differential Operators on Manifolds.

The norms we have given to define the spaces $H_{s}$ depend upon the Fourier transform. In order to get a more invariant definition which can be used to extend these notions to manifolds, we must consider elliptic pseudodifferential operators.

Let $p \in S^{d}(U)$ be a square matrix and let $U_{1}$ be an open set with $\bar{U}_{1} \subset U$. We say that $p$ is elliptic on $U_{1}$ if there exists an open subset $U_{2}$ with $\bar{U}_{1} \subset U_{2} \subset \bar{U}_{2} \subset U$ and if there exists $q \in S^{-d}$ such that $p q-I \in S^{-\infty}$ and $q p-I \in S^{-\infty}$ over $U_{2}$. (To say that $r \in S^{-\infty}$ over $U_{2}$ simply means the estimates of section 1.2 hold over $U_{2}$. Equivalently, we assume $\phi r \in S^{-\infty}$ for every $\left.\phi \in C_{0}^{\infty}\left(U_{2}\right)\right)$. This constant technical fuss over domains will be eliminated very shortly when we pass to considering compact manifolds; the role of $U_{2}$ is to ensure uniform estimates over $U_{1}$.

It is clear that $p$ is elliptic over $U_{1}$ if and only if there exists constants $C_{0}$ and $C_{1}$ such that $p(x, \xi)$ is invertible for $|\xi| \geq C_{0}$ and

$$
\left|p(x, \xi)^{-1}\right| \leq C_{1}(1+|\xi|)^{-d} \quad \text { for }|\xi| \geq C_{0}, \quad x \in U_{2}
$$

We define $q=\phi(\xi) p^{-1}(x, \xi)$ where $\phi(\xi)$ is a cut-off function identically 0 near $\xi=0$ and identically 1 near $\xi=\infty$. We used similar cutoff functions in the proof of Lemma 1.2.8. Furthermore, if $p_{0} \in S^{d-1}$, then $p$ is elliptic if and only if $p+p_{0}$ is elliptic; adding lower order terms does not alter the ellipticity. If $p$ is a polynomial and $P$ is a differential operator, then $p$ is elliptic if and only if the leading symbol $\sigma_{L}(p)=\sum_{|\alpha|=d} p_{\alpha}(x) \xi^{\alpha}$ is invertible for $\xi \neq 0$.

There exist elliptic operators of all orders. Let $\phi(x) \in C_{0}^{\infty}$ and define the symbol $p(x, \xi)=\phi(x)\left(1+|\xi|^{2}\right)^{d / 2} I$, then this is an elliptic symbol of order $d$ whenever $\phi(x) \neq 0$.

Lemma 1.3.1. Let $P \in \Psi^{d}(U)$ be elliptic over $U_{1}$ then:
(a) There exists $Q \in \Psi^{-d}(U)$ such that $P Q-I \sim 0$ and $Q P-I \sim 0$ over $U_{1}$ (i.e., $\phi(P Q-I)$ and $\phi(Q P-I)$ are infinitely smoothing for any $\left.\phi \in C_{0}^{\infty}\left(U_{2}\right)\right)$.
(b) $P$ is hypo-elliptic over $U_{1}$, i.e., if $f \in H_{s}$ and if $P f$ is smooth over $U_{1}$ then $f$ is smooth over $U_{1}$.
(c) There exists a constant $C$ such that $|f|_{d} \leq C\left(|f|_{0}+|P f|_{0}\right)$ for $f \in$ $C_{0}^{\infty}\left(U_{1}\right)$. (Gärding's inequality).

Proof: We will define $Q$ to have symbol $q_{0}+q_{1}+\cdots$ where $q_{j} \in S^{-d-j}$. We try to solve the equation

$$
\sigma(P Q-I) \sim \sum_{\alpha, j} d_{\xi}^{\alpha} p \cdot D_{x}^{\alpha} q_{j} / \alpha!-I \sim 0
$$

When we decompose this sum into elements of $S^{-k}$, we conclude we must solve

$$
\sum_{|\alpha|+j=k} d_{\xi}^{\alpha} p \cdot D_{x}^{\alpha} q_{j} / \alpha!= \begin{cases}I & \text { if } k=0 \\ 0 & \text { if } k \neq 0\end{cases}
$$

We define $q_{0}=q$ and then solve the equation inductively to define:

$$
q_{k}=-q \cdot \sum_{\substack{|\alpha|+j=k \\ j<k}} d_{\xi}^{\alpha} p \cdot D_{x}^{\alpha} q_{j} / \alpha!.
$$

This defines $Q$ so $\sigma(P Q-I) \sim 0$ over $U_{2}$. Similarly we could solve $\sigma\left(Q_{1} P-\right.$ $I) \sim 0$ over $U_{2}$. We now compute $\sigma\left(Q-Q_{1}\right)=\sigma\left(Q-Q_{1} P Q\right)+\sigma\left(Q_{1} P Q-\right.$ $\left.Q_{1}\right)=\sigma\left(\left(I-Q_{1} P\right) Q\right)+\sigma\left(Q_{1}(Q P-I)\right) \sim 0$ over $U_{2}$ so that in fact $Q$ and $Q_{1}$ agree modulo infinitely smoothing operators. This proves (a).

Let $f \in H_{s}$ with $P f$ smooth over $U_{1}$, and choose $\phi \in C_{0}\left(U_{1}\right)$. We compute:

$$
\phi f=\phi(I-Q P) f+\phi Q P f
$$

As $\phi(I-Q P) \sim 0, \phi(I-Q P) f$ is smooth. Since $P f$ is smooth over $U_{1}$, $\phi Q P f$ is smooth since $Q$ is pseudo-local. Thus $\phi f$ is smooth which proves (b).

Finally, we choose $\phi \in C_{0}^{\infty}\left(U_{2}\right)$ to be identically 1 on $U_{1}$. Then if $f \in C_{0}^{\infty}\left(U_{1}\right)$,

$$
|f|_{d}=|\phi f|_{d}=|\phi(I-Q P) f+\phi Q P f|_{d} \leq|\phi(I-Q P) f|_{d}+|\phi Q P f|_{d}
$$

We estimate the first norm by $C|f|_{0}$ since $\phi(I-Q P)$ is an infinitely smoothing operator. We estimate the second norm by $C|P f|_{0}$ since $\phi Q$ is a bounded map from $L^{2}$ to $H^{d}$. This completes the proof.

We note (c) is immediate if $d<0$ since $\left.\right|_{d} \leq\left.\right|_{0}$. If $d>0,|f|_{0}+|P f|_{0} \leq$ $C\left(|f|_{d}\right)$ so this gives a equivalent norm on $H^{d}$.

We now consider the effect of changes of coordinates on our class of pseudo-differential operators. Let $h: U \rightarrow \widetilde{U}$ be a diffeomorphism. We define $h^{*}: C^{\infty}(\widetilde{U}) \rightarrow C^{\infty}(U)$ by $h^{*} f(x)=f(h(x))$. If $P$ is a linear operator on $C^{\infty}(U)$, we define $h_{*} P$ acting on $C^{\infty}(\widetilde{U})$ by $\left(h_{*} P\right) f=\left(h^{-1}\right)^{*} P\left(h^{*} f\right)$. The fundamental lemma we shall need is the following:
Lemma 1.3.2. Let $h: U \rightarrow \widetilde{U}$ be a diffeomorphism. Then:
(a) If $P \in \Psi^{d}(U)$ then $h_{*} P \in \Psi^{d}(\widetilde{U})$. Let $p=\sigma(P)$ and define $h(x)=x_{1}$ and $d h(x)^{t} \xi_{1}=\xi$. Let $p_{1}\left(x_{1}, \xi_{1}\right)=p(x, \xi)$ then $\sigma\left(h_{*} P\right)-p_{1} \in S^{d-1}$.
(b) Let $U_{1}$ be an open subset with $\bar{U}_{1} \subset U$. There exists a constant $C$ such that $\left|h^{*} f\right|_{d} \leq C|f|_{d}$ for all $f \in C_{0}^{\infty}\left(h\left(U_{1}\right)\right)$. In other words, the Sobolev spaces are invariant.
Proof: The first step is to localize the problem. Let $\left\{\phi_{i}\right\}$ be a partition of unity and let $P_{i j}=\phi_{i} P \phi_{j}$ so $P=\sum_{i, j} P_{i j}$. If the support of $\phi_{i}$ is
disjoint from the support of $\phi_{j}$, then $P_{i j}$ is an infinitely smoothing operator with a smooth kernel $K_{i j}(x, y)$ by Lemma 1.2.6. Therefore $h_{*} P_{i j}$ is also given by a smooth kernel and is a pseudo-differential operator by Lemma 1.2.4. Consequently, we may restrict attention to pairs $(i, j)$ such that the supports of $\phi_{i}$ and $\phi_{j}$ intersect. We assume henceforth $P$ is defined by a symbol $p(x, \xi, y)$ where $p$ has arbitrarily small support in $(x, y)$.

We first suppose $h$ is linear to motivate the constructions of the general case. Let $h(x)=h x$ where $h$ is a constant matrix. We equate:

$$
h x=x_{1}, \quad h y=y_{1}, \quad h^{t} \xi_{1}=\xi
$$

and define

$$
p_{1}\left(x_{1}, \xi_{1}, y_{1}\right)=p(x, \xi, y)
$$

(In the above, $h^{t}$ denotes the matrix transpose of $h$ ). If $f \in C_{0}^{\infty}(\widetilde{U})$, we compute:

$$
\begin{aligned}
\left(h_{*} P\right) f\left(x_{1}\right)= & \int e^{i(x-y) \cdot \xi} p(x, \xi, y) f(h y) d y d \xi \\
= & \int e^{i h^{-1}\left(x_{1}-y_{1}\right) \cdot \xi} p\left(h^{-1} x_{1}, \xi, h^{-1} y_{1}\right) f\left(y_{1}\right) \\
& \quad \times|\operatorname{det}(h)|^{-1} d y_{1} d \xi
\end{aligned}
$$

We now use the identities $h^{-1}\left(x_{1}-y_{1}\right) \cdot \xi=\left(x_{1}-y_{1}\right) \cdot \xi_{1}$ and $|\operatorname{det}(h)| d \xi_{1}=$ $d \xi$ to write:

$$
\begin{aligned}
\left(h_{*} P\right) f\left(x_{1}\right) & =\int e^{i\left(x_{1}-y_{1}\right) \cdot \xi_{1}} p\left(h^{-1} x_{1}, h^{t} \xi_{1}, h^{-1} y_{1}\right) f\left(y_{1}\right) d y_{1} d \xi_{1} \\
& =\int e^{i\left(x_{1}-y_{1}\right) \cdot \xi_{1}} p_{1}\left(x_{1}, \xi_{1}, y_{1}\right) f\left(y_{1}\right) d y_{1} d \xi_{1}
\end{aligned}
$$

This proves that $\left(h_{*} P\right)$ is a pseudo-differential operator on $\widetilde{U}$. Since we don't need to localize in this case, we compute directly that

$$
\sigma\left(h_{*} P\right)\left(x_{1}, \xi_{1}\right)=p\left(h^{-1} x_{1}, h^{t} \xi_{1}\right)
$$

We regard $(x, \xi)$ as giving coordinates for $T^{*} M$ when we expand any covector in the form $\sum \xi_{i} d x^{i}$. This is exactly the transformation for the cotangent space so we may regard $\sigma P$ as being invariantly defined on $T^{*} \mathbf{R}^{m}$.

If $h$ is not linear, the situation is somewhat more complicated. Let

$$
\begin{aligned}
x-y & =h^{-1}\left(x_{1}\right)-h^{-1}\left(y_{1}\right)=\int_{0}^{1} \frac{d}{d t}\left\{h^{-1}\left(t x_{1}+(1-t) y_{1}\right)\right\} d t \\
& =\int_{0}^{1} d\left(h^{-1}\right)\left(t x_{1}+(1-t) y_{1}\right) \cdot\left(x_{1}-y_{1}\right) d t=T\left(x_{1}, y_{1}\right)\left(x_{1}-y_{1}\right)
\end{aligned}
$$

where $T\left(x_{1}, y_{1}\right)$ is a square matrix. If $x_{1}=y_{1}$, then $T\left(x_{1}, y_{1}\right)=d\left(h^{-1}\right)$ is invertible since $h$ is a diffeomorphism. We localize using a partition of unity to suppose henceforth the supports are small enough so $T\left(x_{1}, y_{1}\right)$ is invertible for all points of interest.

We set $\xi_{1}=T\left(x_{1}, y_{1}\right)^{t} \xi$ and compute:

$$
\begin{aligned}
\left(h_{*} P\right)(f)\left(x_{1}\right) & =\int e^{i(x-y) \cdot \xi} p(x, \xi, y) f(h y) d y d \xi \\
& =\int e^{i T\left(x_{1}, y_{1}\right)\left(x_{1}-y_{1}\right) \cdot \xi} p\left(h^{-1} x_{1}, \xi, h^{-1} y_{1}\right) f\left(y_{1}\right) J d y_{1} d \xi \\
& =\int e^{i\left(x_{1}-y_{1}\right) \cdot \xi_{1}} p_{1}\left(x_{1}, \xi_{1}, y_{1}\right) f\left(y_{1}\right) \\
& \quad \times J\left|\operatorname{det} T\left(x_{1}, y_{1}\right)\right|^{-1} d y_{1} d \xi_{1}
\end{aligned}
$$

where $J=\left|\operatorname{det}\left(d h^{-1}\right)\right|=\left|\operatorname{det} T\left(y_{1}, y_{1}\right)\right|$. By Lemma 1.2.2, this defines a pseudo-differential operator of order $d$ such that $\sigma\left(h_{*} P\right)=p_{1}$ modulo $S^{d-1}$ which completes the proof of (a). Since $|d h|$ is uniformly bounded on $U_{1}$, $|f|_{0} \leq C\left|h^{*} f\right|_{0}$ and $\left|h^{*} f\right|_{0} \leq C|f|_{0}$. If $P$ is elliptic, then $h_{*} P$ is elliptic of the same order. For $d>0$, choose $P$ elliptic of order $d$ and compute:

$$
\left|h^{*} f\right|_{d} \leq C\left(\left|h^{*} f\right|_{0}+\left|P h^{*} f\right|_{0}\right) \leq C\left(|f|_{0}+\left|\left(h_{*} P\right) f\right|_{0}\right) \leq C|f|_{d}
$$

which completes the proof of (b) if $d \geq 0$. The result for $d \leq 0$ follows by duality using Lemma 1.1.6(b).

We introduce the spaces $S^{d} / S^{d-1}$ and define $\sigma_{L}(P)$ to be the element defined by $\sigma(P)$ in this quotient. Let $P$ and $Q$ be pseudo-differential operators of order $d_{1}$ and $d_{2}$ Then $P Q$ is a pseudo-differential operator of order $d_{1}+d_{2}$ and $\sigma_{L}(P Q)=\sigma_{L}(P) \sigma_{L}(Q)$ since the remaining terms in the asymptotic series are of lower order. Similarly $\sigma_{L}\left(P^{*}\right)=\sigma_{L}(P)^{*}$. If we define $(x, \xi)$ as coordinates for $T^{*}\left(\mathbf{R}^{m}\right)$ by representing a cotangent vector at a point $x$ in the form $\sum \xi_{i} d x^{i}$, then Lemma 1.3.2 implies $\sigma_{L}(P)$ is invariantly defined on $T^{*}\left(\mathbf{R}^{m}\right)$. If $P=\sum_{\alpha} p_{\alpha} D_{x}^{\alpha}$ is a differential operator there is a natural identification of $\sum_{|\alpha|=d} p_{\alpha} \xi^{\alpha}$ with the image of $p$ in $S^{d} / S^{d-1}$ so this definition of the leading symbol agrees with that given earlier.

We now extend the results of section 1.2 to manifolds. Let $M$ be a smooth compact Riemannian manifold without boundary. Let $m$ be the dimension of $M$ and let dvol or sometimes simply $d x$ denote the Riemanian measure on $M$. In Chapter 2, we will use the notation $\mid$ dvol $\mid$ to denote this measure in order to distinguish between measures and $m$-forms, but we shall not bother with this degree of formalism here. We restrict to scalars first. Let $C^{\infty}(M)$ be the space of smooth functions on $M$ and let $P: C^{\infty}(M) \rightarrow C^{\infty}(M)$ be a linear operator. We say that $P$ is a pseudodifferential operator of order $d$ and write $P \in \Psi_{d}(M)$ if for every open chart
$U$ on $M$ and for every $\phi, \psi \in C_{0}^{\infty}(U)$, the localized operator $\phi P \psi \in \Psi_{d}(U)$. We say that $P$ is elliptic if $\phi P \psi$ is elliptic where $\phi \psi(x) \neq 0$. If $Q \in \Psi_{d}(U)$, we let $P=\phi Q \psi$ for $\phi, \psi \in C_{0}^{\infty}(U)$. Lemma 1.3.2 implies $P$ is a pseudodifferential operator on $M$ so there exists operators of all orders on $M$. We define:

$$
\Psi(M)=\bigcup_{d} \Psi_{d}(M) \quad \text { and } \quad \Psi_{-\infty}(M)=\bigcap_{d} \Psi_{d}(M)
$$

to be the set of all pseudo-differential operators on $M$ and the set of infinitely smoothing operators on $M$.

In any coordinate system, we define $\sigma(P)$ to the symbol of the operator $\phi P \phi$ where $\phi=1$ near the point in question; this is unique modulo $S^{-\infty}$. The leading symbol is invariantly defined on $T^{*} M$, but the total symbol changes under the same complicated transformation that the total symbol of a differential operator does under coordinate transformations. Since we shall not need this transformaton law, we omit the statement; it is implicit in the computations performed in Lemma 1.3.2.

We define $L^{2}(M)$ using the $L^{2}$ inner product

$$
(f, g)=\int_{M} f(x) \bar{g}(x) d x, \quad|f|_{0}^{2}=(f, f)
$$

We let $L^{2}(M)$ be the completion of $C^{\infty}(M)$ in this norm. Let $P: C^{\infty}(M) \rightarrow$ $C^{\infty}(M)$. We let $P^{*}$ be defined by $(P f, g)=\left(f, P^{*} g\right)$ if such a $P^{*}$ exists. Lemmas 1.3.2 and 1.2.3 imply that:

Lemma 1.3.3.
(a) If $P \in \Psi_{d}(M)$, then $P^{*} \in \Psi_{d}(M)$ and $\sigma_{L}\left(P^{*}\right)=\sigma_{L}(P)^{*}$. In any coordinate chart, $\sigma\left(P^{*}\right)$ has a asymptotic expansion given by Lemma 1.2.3(a). (b) If $P \in \Psi_{d}(M)$ and $Q \in \Psi_{e}(M)$, then $P Q \in \Psi_{d+e}(M)$ and $\sigma_{L}(P Q)=$ $\sigma_{L}(P) \sigma_{L}(Q)$. In any coordinate chart $\sigma(P Q)$ has an asymptotic expansion given in Lemma 1.2.3(b).

We use a partition of unity to define the Sobolev spaces $H_{s}(M)$. Cover $M$ by a finite number of coordinate charts $U_{i}$ with diffeomorphisms $h_{i}: O_{i} \rightarrow$ $U_{i}$ where the $O_{i}$ are open subsets of $\mathbf{R}^{m}$ with compact closure. If $f \in$ $C_{0}^{\infty}\left(U_{i}\right)$, we define:

$$
|f|_{s}^{(i)}=\left|h_{i}^{*} f\right|_{s}
$$

where we shall use the superscript ${ }^{(i)}$ to denote the localized norm. Let $\left\{\phi_{i}\right\}$ be a partition of unity subordinate to this cover and define:

$$
|f|_{s}=\sum\left|\phi_{i} f\right|_{s}^{(i)}
$$

If $\psi \in C^{\infty}(M)$, we note that $\left|\phi_{i} \psi f\right|_{s}^{(i)} \leq C\left|\phi_{i} f\right|_{s}^{(i)}$ since multiplication by $\psi$ defines a $\Psi \mathrm{DO}$ of order 0 . Suppose $\left\{U_{j}^{\prime}, O_{j}^{\prime}, h_{j}^{\prime}, \phi_{j}^{\prime}\right\}$ is another possible choice to define $\left.\right|_{s} ^{\prime}$. We estimate:

$$
\left|\phi_{j}^{\prime} f\right|_{s}^{(j)} \leq \sum_{i}\left|\phi_{j}^{\prime} \phi_{i} f\right|_{s}^{(j)}
$$

Since $\phi_{j}^{\prime} \phi_{i} f \in C_{0}^{\infty}\left(U_{i} \cap U_{j}^{\prime}\right)$, we can use Lemma 1.3.2 (b) to estimate $\left.\right|_{s} ^{\prime(j)}$ by $\left.\right|_{s} ^{(i)}$ so

$$
\left|\phi_{j}^{\prime} f\right|_{s}^{(j)} \leq \sum_{i} C\left|\phi_{j}^{\prime} \phi_{i} f\right|_{s}^{(i)} \leq C\left|\phi_{i} f\right|_{s}^{(i)} \leq C|f|_{s}
$$

so that $|f|_{s}^{\prime} \leq C|f|_{s}$. Similarly $|f|_{s} \leq C|f|_{s}^{\prime}$. This shows these two norms are equivalent so $H_{s}(M)$ is defined independent of the choices made.

We note that:

$$
\sum_{i}\left\{\left|\phi_{i} f\right|_{s}^{(i)}\right\}^{2} \leq|f|_{s}^{2} \leq C \sum_{i}\left\{\left|\phi_{i} f\right|_{s}^{(i)}\right\}^{2}
$$

so by using this equivalent norm we conclude the $H_{s}(M)$ are topologically Hilbert spaces. If $\left\{\psi_{i}\right\}$ are given subordinate to the cover $U_{i}$ with $\psi_{i} \geq 0$ and $\psi=\sum_{i} \psi_{i}>0$, we let $\phi_{i}=\psi_{i} / \psi$ and compute:

$$
\begin{aligned}
\sum_{i}\left|\psi_{i} f\right|_{s}^{(i)} & =\sum_{i}\left|\psi \phi_{i} f\right|_{s}^{(i)} \leq C \sum_{i}\left|\phi_{i} f\right|_{s}^{(i)}=C|f|_{s} \\
\sum_{i}\left|\phi_{i} f\right|_{s}^{(i)} & =\sum_{i}\left|\psi^{-1} \psi_{i} f\right|_{s}^{(i)} \leq C \sum_{i}\left|\psi_{i} f\right|_{s}^{(i)}
\end{aligned}
$$

to see the norm defined by $\sum_{i}\left|\psi_{i} f\right|_{s}^{(i)}$ is equivalent to the norm $|f|_{s}$ as well.

Lemma 1.1.5 implies that $|f|_{t} \leq|f|_{s}$ if $t<s$ and that the inclusion of $H_{s}(M) \rightarrow H_{t}(M)$ is compact. Lemma 1.1.7 implies given $s>t>u$ and $\varepsilon>0$ we can estimate:

$$
|f|_{t} \leq \varepsilon|f|_{s}+C(\varepsilon)|f|_{u}
$$

We assume the coordinate charts $U_{i}$ are chosen so the union $U_{i} \cup U_{j}$ is also contained in a larger coordinate chart for all $(i, j)$. We decompose $p \in \Psi_{d}(M)$ as $P=\sum_{i, j} P_{i, j}$ for $P_{i, j}=\phi_{i} P \phi_{j}$. By Lemma 1.2.1 we can estimate:

$$
\left|\phi_{i} P \phi_{j} f\right|_{s}^{(i)} \leq C\left|\phi_{j} f\right|_{s+d}^{(j)}
$$

so $|P f|_{s} \leq C|f|_{s+d}$ and $P$ extends to a continuous map from $H_{s+d}(M) \rightarrow$ $H_{s}(M)$ for all $s$.

We define:

$$
|f|_{\infty, k}=\sum_{i}\left|\phi_{i} f\right|_{\infty, k}^{(i)}
$$

as a measure of the sup norm of the $k^{\text {th }}$ derivatives of $f$. This is independent of particular choices made. Lemma 1.1.4 generalizes to:

$$
|f|_{\infty, k} \leq C|f|_{s} \quad \text { for } s>\frac{m}{2}+k
$$

Thus $H_{s}(M)$ is a subset of $C^{k}(M)$ in this situation.
We choose $\psi_{i} \in C^{\infty}\left(U_{i}\right)$ with $\sum_{i} \psi_{i}^{2}=1$ then:

$$
\begin{aligned}
|(f, g)|=\left|\sum_{i}\left(\psi_{i} f, \psi_{i} g\right)\right| & \leq \sum_{i}\left|\left(\psi_{i} f, \psi_{i} g\right)\right| \\
& \leq C \sum_{i}\left|\psi_{i} f\right|_{s}^{(i)}\left|\psi_{i} g\right|_{-s}^{(i)} \leq C|f|_{s}|g|_{-s}
\end{aligned}
$$

Thus the $L^{2}$ inner product gives a continuous map $H_{s}(M) \times H_{-s}(M) \rightarrow \mathbf{C}$.
Lemma 1.3.4.
(a) The natural inclusion $H_{s} \rightarrow H_{t}$ is compact for $s>t$. Furthermore, if $s>t>u$ and if $\varepsilon>0$, then $|f|_{t} \leq \varepsilon|f|_{s}+C(\varepsilon)|f|_{u}$.
(b) If $s>k+\frac{m}{2}$ then $H_{s}(M)$ is contained in $C^{k}(M)$ and we can estimate $|f|_{\infty, k} \leq C|f|_{s}$.
(c) If $P \in \Psi_{d}(M)$ then $P: H_{s+d}(M) \rightarrow H_{s}(M)$ is continuous for all $s$.
(d) The pairing $H_{s}(M) \times H_{-s}(M) \rightarrow \mathbf{C}$ given by the $L^{2}$ inner product is a perfect pairing.

Proof: We have proved every assertion except the fact (d) that the pairing is a perfect pairing. We postpone this proof briefly until after we have discussed elliptic $\Psi D O$ 's.

The sum of two elliptic operators need not be elliptic. However, the sum of two elliptic operators with positive symbols is elliptic. Let $P_{i}$ have symbol $\left(1+\left|\xi_{i}\right|^{2}\right)^{d / 2}$ on $U_{i}$ and let $\phi_{i}$ be a partition of unity. $P=\sum_{i} \phi_{i} P_{i} \phi_{i}$; this is an elliptic $\Psi \mathrm{DO}$ of order $d$ for any $d$, so elliptic operators exist. We let $P$ be an elliptic $\Psi \mathrm{DO}$ of order $d$ and let $\psi_{i}$ be identically 1 on the support of $\phi_{i}$. We use these functions to construct $Q \in \Psi_{-d}(M)$ so $P Q-I \in \Psi_{-\infty}(M)$ and $Q P-I \in \Psi_{-\infty}(M)$. In each coordinate chart, let $P_{i}=\psi_{i} P \psi_{i}$ then $P_{i}-P_{j} \in \Psi_{-\infty}$ on the support of $\phi_{i} \phi_{j}$. We construct $Q_{i}$ as the formal inverse to $P_{i}$ on the support of $\phi_{i}$, then $Q_{i}-Q_{j} \in \Psi_{-\infty}$ on the support of $\phi_{i} \phi_{j}$ since the formal inverse is unique. Modulo $\Psi_{-\infty}$ we have $P=\sum_{i} \phi_{i} P \sim \sum_{i} \phi_{i} P_{i}$. We define $Q=\sum_{j} Q_{j} \phi_{j}$ and note $Q$ has the desired properties.

It is worth noting we could also construct the formal inverse using a Neumann series. By hypothesis, there exists $q$ so $q p-I \in S^{-1}$ and $p q-I \in$
$S^{-1}$, where $p=\sigma_{L}(P)$. We construct $Q_{1}$ using a partition of unity so $\sigma_{L}\left(Q_{1}\right)=q$. We let $Q^{r}$ and $Q^{l}$ be defined by the formal series:

$$
\begin{aligned}
& Q^{l}=Q_{1}\left\{\sum_{k}(-1)^{k}\left(P Q_{1}-I\right)^{k}\right\} \\
& Q^{r}=\left\{\sum_{k}(-1)^{k}\left(Q_{1} P-I\right)^{k}\right\} Q_{1}
\end{aligned}
$$

to construct formal left and right inverses so $Q=Q^{r}=Q^{l}$ modulo $\Psi_{-\infty}$.
Let $P$ be elliptic of order $d>0$. We estimate:

$$
|f|_{d} \leq|(Q P-I) f|_{d}+|Q P f|_{d} \leq C|f|_{0}+C|P f|_{0} \leq C|f|_{d}
$$

so we could define $H_{d}$ using the norm $|f|_{d}=|f|_{0}+|P f|_{0}$. We specialize to the following case. Let $Q$ be elliptic of order $d / 2$ and let $P=Q^{*} Q+1$. Then $Q^{*} Q$ is self-adjoint and non-negative so we can estimate $|f|_{0} \leq|P f|_{0}$ and we can define $|f|_{d}=|P f|_{0}$ in this case. Consequently:

$$
|f|_{d}^{2}=(P f, P f)=\left(f, P^{*} P f\right)=(f, g),
$$

for $g=P^{*} P f$. Since $\left|P^{*} P f\right|_{-d} \leq C|f|_{d}$, we conclude:

$$
|f|_{d}=(P f, P f) /|f|_{d} \leq C(f, g) /|g|_{-d} \leq C \sup _{h}|(f, h)| /|h|_{-d} .
$$

Since the pairing of $H_{d}$ with $H_{-d}$ is continuous, this proves $H_{d}=H_{-d}^{*}$. Topologically these are Hilbert spaces so we see dually that $H_{d}^{*}=H_{-d}$. This completes the proof of Lemma 1.3.4. We also note that we have proved:

Lemma 1.3.5. Let $P \in \Psi_{d}$ be elliptic. Then there exists $Q \in \Psi_{-d}$ so $P Q-I \in \Psi_{-\infty}$ and $Q P-I \in \Psi_{-\infty} . \quad P$ is hypoelliptic. If $d>0$, we can define $h_{d}$ by using the norm $|f|_{0}+|P f|_{0}$ and define $H_{-d}$ by duality.

If $V$ is a vector bundle, we cover $M$ by coordinate charts $U_{i}$ over which $V$ is trivial. We use this cover to define $H_{s}(V)$ using a partition of unity. We shall always assume $V$ has a given fiber metric so $L^{2}(V)$ is invariantly defined. $P: C^{\infty}(V) \rightarrow C^{\infty}(W)$ is a $\Psi \mathrm{DO}$ of order $d$ if $\phi P \psi$ is given by a matrix of $d^{\text {th }}$ order $\Psi$ DO's for $\phi, \psi \in C_{0}^{\infty}(U)$ for any coordinate chart $U$ over which $V$ and $W$ are trivial. Lemmas 1.3.4 and 1.3.5 generalize immediately to this situation.

### 1.4. Fredholm Operators and the Index of a Fredholm Operator.

Elliptic $\Psi$ DO's are invertible modulo $\Psi_{-\infty}$. Lemma 1.3.4 will imply that elliptic $\Psi D O$ 's are invertible modulo compact operators and that such operators are Fredholm. We briefly review the facts we shall need concerning Fredholm and compact operators.

Let $H$ be a Hilbert space and let $\operatorname{END}(H)$ denote the space of all bounded linear maps $T: H \rightarrow H$. There is a natural norm on $\operatorname{END}(H)$ defined by:

$$
|T|=\sup _{x \in H} \frac{|T x|}{|x|}
$$

where the sup ranges over $x \neq 0$. $\mathrm{END}(H)$ becomes a Banach space under this norm. The operations of addition, composition, and taking adjoint are continuous. We let GL $(H)$ be the subset of $\operatorname{END}(H)$ consisting of maps $T$ which are 1-1 and onto. The inverse boundedness theorem shows that if $T \in \operatorname{END}(H)$ is $1-1$ and onto, then there exists $\varepsilon>0$ such that $|T x| \geq \varepsilon|x|$ so $T^{-1}$ is bounded as well. The Neuman series:

$$
(1-z)^{-1}=\sum_{k=0}^{\infty} z^{k}
$$

converges for $|z|<1$. If $|I-T|<1$, we may express $T=I-(I-T)$. If we define

$$
S=\sum_{k=0}^{\infty}(I-T)^{k}
$$

then this converges in $\operatorname{END}(H)$ to define an element $S \in \operatorname{END}(H)$ so $S T=T S=I$. Furthermore, this shows $|T|^{-1} \leq(1-|I-T|)^{-1}$ so GL $(H)$ contains an open neighborhood of $I$ and the map $T \rightarrow T^{-1}$ is continuous there. Using the group operation on GL $(H)$, we see that $\mathrm{GL}(H)$ is a open subset of $\operatorname{END}(H)$ and is a topological group.

We say that $T \in \operatorname{END}(H)$ is compact if $T$ maps bounded sets to precompact sets-i.e., if $\left|x_{n}\right| \leq C$ is a bounded sequence, then there exists a subsequence $x_{n_{k}}$ so $T x_{n_{k}} \rightarrow y$ for some $y \in H$. We let $\operatorname{COM}(H)$ denote the set of all compact maps.

Lemma 1.4.1. $\operatorname{COM}(H)$ is a closed 2-sided $*$-ideal of $\operatorname{END}(H)$.
Proof: It is clear the sum of two compact operators is compact. Let $T \in \operatorname{END}(H)$ and let $C \in \operatorname{COM}(H)$. Let $\left\{x_{n}\right\}$ be a bounded sequence in $H$ then $\left\{T x_{n}\right\}$ is also a bounded sequence. By passing to a subsequence, we may assume $C x_{n} \rightarrow y$ and $C T x_{n} \rightarrow z$. Since $T C x_{n} \rightarrow T y$, this implies $C T$ and $T C$ are compact so $\operatorname{COM}(H)$ is a ideal. Next let $C_{n} \rightarrow C$ in $\operatorname{END}(H)$, and let $x_{n}$ be a bounded sequence in $H$. Choose a subsequence
$x_{n}^{1}$ so $C_{1} x_{n}^{1} \rightarrow y^{1}$. We choose a subsequence of the $x_{n}^{1}$ so $C_{2} x_{n}^{2} \rightarrow y^{2}$. By continuing in this way and then using the diagonal subsequence, we can find a subsequence we denote by $x_{n}^{n}$ so $C_{k}\left(x_{n}^{n}\right) \rightarrow y^{k}$ for all $k$. We note $\left|C x_{n}^{n}-C_{k} x_{n}^{n}\right| \leq\left|C-C_{k}\right| c$. Since $\left|C-C_{k}\right| \rightarrow 0$ this shows the sequence $C x_{n}^{n}$ is Cauchy so $C$ is compact and $\operatorname{COM}(H)$ is closed. Finally let $C \in \operatorname{COM}(H)$ and suppose $C^{*} \notin \operatorname{COM}(H)$. We choose $\left|x_{n}\right| \leq 1$ so $\left|C^{*} x_{n}-C^{*} x_{m}\right| \geq \varepsilon>0$ for all $n, m$. We let $y_{n}=C^{*} x_{n}$ be a bounded sequence, then $\left(C y_{n}-C y_{m}, x_{n}-x_{m}\right)=\left|C^{*} x_{n}-C^{*} x_{m}\right|^{2} \geq \varepsilon^{2}$. Therefore $\varepsilon^{2} \leq\left|C y_{n}-C y_{m}\right|\left|x_{n}-x_{m}\right| \leq 2\left|C y_{n}-C y_{m}\right|$ so $C y_{n}$ has no convergent subsequence. This contradicts the assumption $C \in \operatorname{COM}(H)$ and proves $C^{*} \in \operatorname{COM}(H)$.

We shall assume henceforth that $H$ is a separable infinite dimensional space. Although any two such Hilbert spaces are isomorphic, it is convenient to separate the domain and range. If $E$ and $F$ are Hilbert spaces, we define $\operatorname{HOM}(E, F)$ to be the Banach space of bounded linear maps from $E$ to $F$ with the operator norm. We let $\operatorname{ISO}(E, F)$ be the set of invertible maps in $\operatorname{HOM}(E, F)$ and let $\operatorname{COM}(E, F)$ be the closed subspace of $\operatorname{HOM}(E, F)$ of compact maps. If we choose a fixed isomorphism of $E$ with $F$, we may identify $\operatorname{HOM}(E, F)=\operatorname{END}(E), \operatorname{ISO}(E, F)=\mathrm{GL}(E)$, and $\operatorname{COM}(E, F)=\operatorname{COM}(E) . \quad \mathrm{ISO}(E, F)$ is a open subset of $\operatorname{HOM}(E, F)$ and the operation of taking the inverse is a continuous map from $\operatorname{ISO}(E, F)$ to $\operatorname{ISO}(F, E)$. If $T \in \operatorname{HOM}(E, F)$, we define:

$$
\begin{array}{ll}
\mathrm{N}(T)=\{e \in E: T(E)=0\} & \text { (the null space) } \\
\mathrm{R}(T)=\{f \in F: f=T(e) \text { for some } e \in E\} & \text { (the range). }
\end{array}
$$

$\mathrm{N}(T)$ is always closed, but $\mathrm{R}(T)$ need not be. If $\perp$ denotes the operation of taking orthogonal complement, then $\mathrm{R}(T)^{\perp}=\mathrm{N}\left(T^{*}\right)$. We let $\operatorname{FRED}(E, F)$ be the subset of $\operatorname{HOM}(E, F)$ consisting of operators invertible modulo compact operators:

$$
\begin{aligned}
\operatorname{FRED}(E, F)=\{T \in & \operatorname{HOM}(E, F): \exists S_{1}, S_{2} \in \operatorname{HOM}(F, E) \text { so } \\
& \left.S_{1} T-I \in \operatorname{COM}(E) \text { and } T S_{2}-I \in \operatorname{COM}(F)\right\} .
\end{aligned}
$$

We note this condition implies $S_{1}-S_{2} \in \operatorname{COM}(F, E)$ so we can assume $S_{1}=$ $S_{2}$ if we like. The following lemma provides another useful characterization of $\operatorname{FRED}(E, F)$ :
Lemma 1.4.2. The following are equivalent:
(a) $T \in \operatorname{FRED}(E, F)$;
(b) $T \in \operatorname{END}(E, F)$ has $\operatorname{dim} \mathrm{N}(T)<\infty$, $\operatorname{dim} \mathrm{N}\left(T^{*}\right)<\infty, \mathrm{R}(T)$ is closed, and $\mathrm{R}\left(T^{*}\right)$ is closed.

Proof: Let $T \in \operatorname{FRED}(E, F)$ and let $x_{n} \in \mathrm{~N}(T)$ with $\left|x_{n}\right|=1$. Then $x_{n}=\left(I-S_{1} T\right) x_{n}=C x_{n}$. Since $C$ is compact, there exists a convegent
subsequence. This implies the unit sphere in $\mathrm{N}(T)$ is compact so $\mathrm{N}(T)$ is finite dimensional. Next let $y_{n}=T x_{n}$ and $y_{n} \rightarrow y$. We may assume without loss of generality that $x_{n} \in \mathrm{~N}(T)^{\perp}$. Suppose there exists a constant $C$ so $\left|x_{n}\right| \leq C$. We have $x_{n}=S_{1} y_{n}+\left(I-S_{1} T\right) x_{n}$. Since $S_{1} y_{n} \rightarrow S_{1} y$ and since ( $I-S_{1} T$ ) is compact, we can find a convergent subsequence so $x_{n} \rightarrow x$ and hence $y=\lim _{n} y_{n}=\lim _{n} T x_{n}=T x$ is in the range of $T$ so $\mathrm{R}(T)$ will be closed. Suppose instead $\left|x_{n}\right| \rightarrow \infty$. If $x_{n}^{\prime}=x_{n} /\left|x_{n}\right|$ we have $T x_{n}^{\prime}=y_{n} /\left|x_{n}\right| \rightarrow 0$. We apply the same argument to find a subsequence $x_{n}^{\prime} \rightarrow x$ with $T x=0,|x|=1$, and $x \in \mathrm{~N}(T)^{\perp}$. Since this is impossible, we conclude $\mathrm{R}(T)$ is closed (by passing to a subsequence, one of these two possibilities must hold). Since $T^{*} S_{1}^{*}-I=C_{1}^{*}$ and $S_{2}^{*} T^{*}-I=C_{2}^{*}$ we conclude $T^{*} \in \operatorname{FRED}(F, E)$ so $\mathrm{N}\left(T^{*}\right)$ is finite dimensional and $\mathrm{R}\left(T^{*}\right)$ is closed. This proves (a) implies (b).

Conversely, suppose $\mathrm{N}(T)$ and $\mathrm{N}\left(T^{*}\right)$ are finite dimensional and that $\mathrm{R}(T)$ is closed. We decompose:

$$
E=\mathrm{N}(T) \oplus \mathrm{N}(T)^{\perp} \quad F=\mathrm{N}\left(T^{*}\right) \oplus \mathrm{R}(T)
$$

where $T: \mathrm{N}(T)^{\perp} \rightarrow \mathrm{R}(T)$ is 1-1 and onto. Consequently, we can find a bounded linear operator $S$ so $S T=I$ on $\mathrm{N}(T)^{\perp}$ and $T S=I$ on $\mathrm{R}(T)$. We extend $S$ to be zero on $\mathrm{N}\left(T^{*}\right)$ and compute

$$
S T-I=\pi_{\mathrm{N}(T)} \quad \text { and } \quad T S-I=\pi_{\mathrm{N}\left(T^{*}\right)}
$$

where $\pi$ denotes orthogonal projection on the indicated subspace. Since these two projections have finite dimensional range, they are compact which proves $T \in \operatorname{FRED}(E, F)$.

If $T \in \operatorname{FRED}(E, F)$ we shall say that $T$ is Fredholm. There is a natural law of composition:

Lemma 1.4.3.
(a) If $T \in \operatorname{FRED}(E, F)$ then $T^{*} \in \operatorname{FRED}(F, E)$.
(b) If $T_{1} \in \operatorname{FRED}(E, F)$ and $T_{2} \in \operatorname{FRED}(F, G)$ then $T_{2} T_{1} \in \operatorname{FRED}(E, G)$.

Proof: (a) follows from Lemma 1.4.2. If $S_{1} T_{1}-I \in C(E)$ and $S_{2} T_{2}-I \in$ $C(F)$ then $S_{1} S_{2} T_{2} T_{1}-I=S_{1}\left(S_{2} T_{2}-I\right) T_{1}+\left(S_{1} T_{1}-I\right) \in C(E)$. Similarly $T_{2} T_{1} S_{1} S_{2}-I \in C(E)$.

If $T \in \operatorname{FRED}(E, F)$, then we define:

$$
\operatorname{index}(T)=\operatorname{dim} \mathrm{N}(T)-\operatorname{dim} \mathrm{N}\left(T^{*}\right)
$$

We note that $\operatorname{ISO}(E, F)$ is contained in $\operatorname{FRED}(E, F)$ and that $\operatorname{index}(T)=0$ if $T \in \operatorname{ISO}(E, F)$.

Lemma 1.4.4.
(a) $\operatorname{index}(T)=-\operatorname{index}\left(T^{*}\right)$.
(b) If $T \in \operatorname{FRED}(E, F)$ and $S \in \operatorname{FRED}(F, G)$ then

$$
\operatorname{index}(S T)=\operatorname{index}(T)+\operatorname{index}(S)
$$

(c) $\operatorname{FRED}(E, F)$ is a open subset of $\operatorname{HOM}(E, F)$.
(d) index: $\operatorname{FRED}(E, F) \rightarrow \mathbf{Z}$ is continuous and locally constant.

Proof: (a) is immediate from the definition. We compute

$$
\begin{aligned}
\mathrm{N}(S T) & =\mathrm{N}(T) \oplus T^{-1}(\mathrm{R}(T) \cap \mathrm{N}(S)) \\
\mathrm{N}\left(T^{*} S^{*}\right) & =\mathrm{N}\left(S^{*}\right) \oplus\left(S^{*}\right)^{-1}\left(\mathrm{R}\left(S^{*}\right) \cap \mathrm{N}\left(T^{*}\right)\right) \\
& =\mathrm{N}\left(S^{*}\right) \oplus\left(S^{*}\right)^{-1}\left(\mathrm{R}(T)^{\perp} \cap \mathrm{N}(S)^{\perp}\right)
\end{aligned}
$$

so that

$$
\begin{aligned}
\operatorname{index}(S T)= & \operatorname{dim} \mathrm{N}(T)+\operatorname{dim}(\mathrm{R}(T) \cap \mathrm{N}(S))-\operatorname{dim} \mathrm{N}\left(S^{*}\right) \\
& \quad-\operatorname{dim}\left(\mathrm{R}(T)^{\perp} \cap \mathrm{N}(S)^{\perp}\right) \\
= & \operatorname{dim} \mathrm{N}(T)+\operatorname{dim}(\mathrm{R}(T) \cap \mathrm{N}(S))+\operatorname{dim}\left(\mathrm{R}(T)^{\perp} \cap \mathrm{N}(S)\right) \\
& \quad-\operatorname{dim} \mathrm{N}\left(S^{*}\right)-\operatorname{dim}\left(\mathrm{R}(T)^{\perp} \cap \mathrm{N}(S)^{\perp}\right) \\
& \quad-\operatorname{dim}\left(\mathrm{R}(T)^{\perp} \cap \mathrm{N}(S)\right) \\
= & \operatorname{dim} \mathrm{N}(T)+\operatorname{dim} \mathrm{N}(S)-\operatorname{dim} \mathrm{N}\left(S^{*}\right)-\operatorname{dim}\left(\mathrm{R}(T)^{\perp}\right) \\
= & \operatorname{dim} \mathrm{N}(T)+\operatorname{dim} \mathrm{N}(S)-\operatorname{dim} \mathrm{N}\left(S^{*}\right)-\operatorname{dim} \mathrm{N}\left(T^{*}\right) \\
= & \operatorname{index}(S)+\operatorname{index}(T)
\end{aligned}
$$

We prove (c) and (d) as follows. Fix $T \in \operatorname{FRED}(E, F)$. We decompose:

$$
E=\mathrm{N}(T) \oplus \mathrm{N}(T)^{\perp} \quad \text { and } \quad F=\mathrm{N}\left(T^{*}\right) \oplus \mathrm{R}(T)
$$

where $T: \mathrm{N}(T)^{\perp} \rightarrow \mathrm{R}(T)$ is 1-1 onto. We let $\pi_{1}: E \rightarrow \mathrm{~N}(T)$ be orthogonal projection. We define $E_{1}=\mathrm{N}\left(T^{*}\right) \oplus E$ and $F_{1}=\mathrm{N}(T) \oplus F$ to be Hilbert spaces by requiring the decompositon to be orthogonal. Let $S \in \operatorname{HOM}(E, F)$, we define $S_{1} \in \operatorname{HOM}\left(E_{1}, F_{1}\right)$ by:

$$
S_{1}\left(f_{0} \oplus e\right)=\pi_{1}(e) \oplus\left(f_{0}+S_{1}(e)\right)
$$

It is clear $\left|S_{1}-S_{1}^{\prime}\right|=\left|S-S^{\prime}\right|$ so the map $S \rightarrow S_{1}$ defines a continuous map from $\operatorname{HOM}(E, F) \rightarrow \operatorname{HOM}\left(E_{1}, F_{1}\right)$. Let $i_{1}: E \rightarrow E_{1}$ be the natural inclusion and let $\pi_{2}: F_{1} \rightarrow F$ be the natural projection. Since $\mathrm{N}(T)$ and $\mathrm{N}\left(T^{*}\right)$ are finite dimensional, these are Fredholm maps. It is immediate from the definition that if $S \in \operatorname{HOM}(E, F)$ then

$$
S=\pi_{2} S_{1} i_{1}
$$

If we let $T=S$ and decompose $e=e_{0} \oplus e_{1}$ and $f=f_{0} \oplus f_{1}$ for $e_{0} \in \mathrm{~N}(T)$ and $f_{0} \in \mathrm{~N}\left(T^{*}\right)$, then:

$$
T_{1}\left(f_{0} \oplus e_{0} \oplus e_{1}\right)=e_{0} \oplus f_{0} \oplus T e_{1}
$$

so that $T_{1} \in \operatorname{ISO}\left(E_{1}, F_{1}\right)$. Since $\operatorname{ISO}\left(E_{1}, F_{1}\right)$ is an open subset, there exists $\varepsilon>0$ so $|T-S|<\varepsilon$ implies $S_{1} \in \operatorname{ISO}\left(E_{1}, F_{1}\right)$. This implies $S_{1}$ is Fredholm so $S=\pi_{2} S_{1} i_{1}$ is Fredholm and $\operatorname{FRED}(E, F)$ is open. Furthermore, we can compute $\operatorname{index}(S)=\operatorname{index}\left(\pi_{2}\right)+\operatorname{index}\left(S_{1}\right)+\operatorname{index}\left(i_{1}\right)=\operatorname{index}\left(\pi_{2}\right)+$ index $\left(i_{1}\right)=\operatorname{dim} \mathrm{N}(T)-\operatorname{dim} \mathrm{N}\left(T^{*}\right)=\operatorname{index}(T)$ which proves index is locally constant and hence continuous. This completes the proof of the lemma.

We present the following example of an operator with index 1. Let $\phi_{n}$ be orthonormal basis for $L^{2}$ as $n \in \mathbf{Z}$. Define the one sided shift

$$
T \phi_{n}= \begin{cases}\phi_{n-1} & \text { if } n>0 \\ 0 & \text { if } n=0 \\ \phi_{n} & \text { if } n<0\end{cases}
$$

then $T$ is surjective so $\mathrm{N}\left(T^{*}\right)=\{0\}$. Since $\mathrm{N}(T)$ is one dimensional, $\operatorname{index}(T)=1$. Therefore $\operatorname{index}\left(T^{n}\right)=n$ and $\operatorname{index}\left(\left(T^{*}\right)^{n}\right)=-n$. This proves index: $\operatorname{FRED}(E, F) \rightarrow \mathbf{Z}$ is surjective. In the next chapter, we will give several examples of differential operators which have non-zero index.

If we specialize to the case $E=F$ then $\operatorname{COM}(E)$ is a closed two-sided ideal of $\operatorname{END}(E)$ so we can pass to the quotient algebra $\operatorname{END}(E) / \operatorname{COM}(E)$. If $\operatorname{GL}(\operatorname{END}(E) / \operatorname{COM}(E))$ denotes the group of invertible elements and if $\pi: \operatorname{END}(E) \rightarrow \operatorname{END}(E) / \operatorname{COM}(E)$ is the natural projection, then $\operatorname{FRED}(E)$ is $\pi^{-1}$ of the invertible elements. If $C$ is compact and $T$ Fredholm, $T+t C$ is Fredholm for any $t$. This implies index $(T)=$ index $(T+t C)$ so we can extend index: $\operatorname{GL}(\operatorname{END}(E) / \operatorname{COM}(E)) \rightarrow \mathbf{Z}$ as a surjective group homomorphism.

Let $P: C^{\infty}(V) \rightarrow C^{\infty}(W)$ be a elliptic $\Psi \mathrm{DO}$ of order $d$. We construct an elliptic $\Psi$ DO $S$ of order $-d$ with $S: C^{\infty}(W) \rightarrow C^{\infty}(V)$ so that $S P-I$ and $P S-I$ are infinitely smoothing operators. Then

$$
P: H_{s}(V) \rightarrow H_{d-s}(W) \quad \text { and } \quad S: H_{d-s}(W) \rightarrow H_{s}(V)
$$

are continuous. Since $S P-I: H_{s}(V) \rightarrow H_{t}(V)$ is continuous for any $t$, it is compact. Similarly $P S-I$ is a compact operator so both $P$ and $S$ are Fredholm. If $f \in \mathrm{~N}(P)$, then $f$ is smooth by Lemma 1.3.5. Consequently, $\mathrm{N}(P)$ and $\mathrm{N}\left(P^{*}\right)$ are independent of the choice of $s$ and index $(P)$ is invariantly defined. Furthermore, if $P_{\tau}$ is a smooth 1-parameter family of such operators, then index $\left(P_{\tau}\right)$ is independent of the parameter $\tau$ by Lemma 1.4.4. In particular, index $(P)$ only depends on the homotopy type of the leading symbol of $P$. In Chapter 3, we will give a topological formula for index $(P)$ in terms of characteristic classes.

We summarize our conclusions about $\operatorname{index}(P)$ in the following

### 1.4. Fredholm Operators

Lemma 1.4.5. Let $P: C^{\infty}(V) \rightarrow C^{\infty}(W)$ be an elliptic $\Psi D O$ of order $d$ over a compact manifold without boundary. Then:
(a) $\mathrm{N}(P)$ is a finite dimensional subset of $C^{\infty}(V)$.
(b) $P: H_{s}(V) \rightarrow H_{s-d}(W)$ is Fredholm. $P$ has closed range. index $(P)$ does not depend on the particular $s$ chosen.
(c) index $(P)$ only depends on the homotopy type of the leading symbol of $P$ within the class of elliptic $\Psi D O$ 's of order $d$.

### 1.5. Elliptic Complexes, The Hodge Decomposition Theorem, And Poincaré Duality.

Let $V$ be a graded vector bundle. $V$ is a collection of vector bundles $\left\{V_{j}\right\}_{j \in \mathbf{Z}}$ such that $V_{j} \neq\{0\}$ for only a finite number of indices $j$. We let $P$ be a graded $\Psi D O$ of order $d . \quad P$ is a collection of $d^{\text {th }}$ order pseudodifferential operators $P_{j}: C^{\infty}\left(V_{j}\right) \rightarrow C^{\infty}\left(V_{j+1}\right)$. We say that $(P, V)$ is a complex if $P_{j+1} P_{j}=0$ and $\sigma_{L} P_{j+1} \sigma_{L} P_{j}=0$ (the condition on the symbol follows from $P^{2}=0$ for differential operators). We say that $(P, V)$ is elliptic if:

$$
\mathrm{N}\left(\sigma_{L} P_{j}\right)(x, \xi)=\mathrm{R}\left(\sigma_{L} P_{j-1}\right)(x, \xi) \quad \text { for } \xi \neq 0
$$

or equivalently if the complex is exact on the symbol level.
We define the cohomology of this complex by:

$$
H^{j}(V, P)=\mathrm{N}\left(P_{j}\right) / \mathrm{R}\left(P_{j-1}\right)
$$

We shall show later in this section that $H^{j}(V, P)$ is finite dimensional if $(P, V)$ is an elliptic complex. We then define

$$
\operatorname{index}(P)=\sum_{j}(-1)^{j} \operatorname{dim} H^{j}(V, P)
$$

as the Euler characteristic of this elliptic complex.
Choose a fixed Hermitian inner product on the fibers of $V$. We use that inner product together with the Riemannian metric on $M$ to define $L^{2}(V)$. We define adjoints with respect to this structure. If $(P, V)$ is an elliptic complex, we construct the associated self-adjoint Laplacian:

$$
\Delta_{j}=\left(P^{*} P\right)_{j}=P_{j}^{*} P_{j}+P_{j-1} P_{j-1}^{*}: C^{\infty}\left(V_{j}\right) \rightarrow C^{\infty}\left(V_{j}\right)
$$

If $p_{j}=\sigma_{L}\left(P_{j}\right)$, then $\sigma_{L}\left(\Delta_{j}\right)=p_{j}^{*} p_{j}+p_{j-1} p_{j-1}^{*}$. We can also express the condition of ellipticity in terms of $\Delta_{j}$ :

Lemma 1.5.1. Let $(P, V)$ be a $d^{\text {th }}$ order partial differential complex.
Then $(P, V)$ is elliptic if and only if $\Delta_{j}$ is an elliptic operator of order $2 d$ for all $j$.

Proof: We suppose that $(P, V)$ is elliptic; we must check $\sigma_{L}\left(\Delta_{j}\right)$ is nonsingular for $\xi \neq 0$. Suppose $\left(p_{j}^{*} p_{j}+p_{j-1} p_{j-1}^{*}\right)(x, \xi) v=0$. If we dot this equation with $v$, we see $p_{j} v \cdot p_{j} v+p_{j-1}^{*} v \cdot p_{j-1}^{*} v=0$ so that $p_{j} v=p_{j-1}^{*} v=$ 0 . Thus $v \in \mathrm{~N}\left(p_{j}\right)$ so $v \in \mathrm{R}\left(p_{j-1}\right)$ so we can write $v=p_{j-1} w$. Since $p_{j-1}^{*} p_{j-1} w=0$, we dot this equation with $w$ to see $p_{j-1} w \cdot p_{j-1} w=0$ so $v=p_{j-1} w=0$ which proves $\Delta_{j}$ is elliptic. Conversely, let $\sigma_{L}\left(\Delta_{j}\right)$ be nonsingular for $\xi \neq 0$. Since $(P, V)$ is a complex, $\mathrm{R}\left(p_{j-1}\right)$ is a subset of $\mathrm{N}\left(p_{j}\right)$.

Conversely, let $v \in \mathrm{~N}\left(p_{j}\right)$. Since $\sigma_{L}\left(\Delta_{j}\right)$ is non-singular, we can express $v=\left(p_{j}^{*} p_{j}+p_{j-1} p_{j-1}^{*}\right) w$. We apply $p_{j}$ to conclude $p_{j} p_{j}^{*} p_{j} w=0$. We dot this equation with $p_{j} w$ to conclude $p_{j} p_{j}^{*} p_{j} w \cdot p_{j} w=p_{j}^{*} p_{j} w \cdot p_{j}^{*} p_{j} w=0$ so $p_{j}^{*} p_{j} w=0$. This implies $v=p_{j-1} p_{j-1}^{*} w \in \mathrm{R}\left(p_{j-1}\right)$ which completes the proof.

We can now prove the following:
Theorem 1.5.2 (Hodge Decomposition Theorem). Let $(P, V)$ be a $d^{\text {th }}$ order $\Psi D O$ elliptic complex. Then
(a) We can decompose $L^{2}\left(V_{j}\right)=\mathrm{N}\left(\Delta_{j}\right) \oplus \mathrm{R}\left(P_{j-1}\right) \oplus \mathrm{R}\left(P_{j}^{*}\right)$ as an orthogonal direct sum.
(b) $\mathrm{N}\left(\Delta_{j}\right)$ is a finite dimensional vector space and there is a natural isomorphism of $H^{j}(P, V) \simeq \mathrm{N}\left(\Delta_{j}\right)$. The elements of $\mathrm{N}\left(\Delta_{j}\right)$ are smooth sections to $V_{j}$.
Proof: We regard $\Delta_{j}: H_{2 d}\left(V_{j}\right) \rightarrow L^{2}\left(V_{j}\right)$. Since this is elliptic, it is Fredholm. This proves $\mathrm{N}\left(\Delta_{j}\right)$ is finite dimensional. Since $\Delta_{j}$ is hypoelliptic, $\mathrm{N}\left(\Delta_{j}\right)$ consists of smooth sections to $V$. Since $\Delta_{j}$ is self adjoint and Fredholm, $\mathrm{R}\left(\Delta_{j}\right)$ is closed so we may decompose $L^{2}\left(V_{j}\right)=\mathrm{N}\left(\Delta_{j}\right) \oplus \mathrm{R}\left(\Delta_{j}\right)$. It is clear $\mathrm{R}\left(\Delta_{j}\right)$ is contained in the span of $\mathrm{R}\left(P_{j-1}\right)$ and $\mathrm{R}\left(P_{j}^{*}\right)$. We compute the $L^{2}$ inner product:

$$
\left(P_{j}^{*} f, P_{j-1} g\right)=\left(f, P_{j} P_{j-1} g\right)=0
$$

since $P_{j} P_{j-1}=0$. This implies $\mathrm{R}\left(P_{j-1}\right)$ and $\mathrm{R}\left(P_{j}^{*}\right)$ are orthogonal. Let $f \in \mathrm{~N}\left(\Delta_{j}\right)$, we take the $L^{2}$ inner product with $f$ to conclude

$$
0=\left(\Delta_{j} f, f\right)=\left(P_{j} f, P_{j} f\right)+\left(P_{j-1}^{*} f, P_{j-1}^{*} f\right)
$$

so $\mathrm{N}\left(\Delta_{j}\right)=\mathrm{N}\left(P_{j}\right) \cap \mathrm{N}\left(P_{j-1}^{*}\right)$. This implies $\mathrm{R}\left(\Delta_{j}\right)$ contains the span of $\mathrm{R}\left(P_{j-1}\right)$ and $\mathrm{R}\left(P_{j}^{*}\right)$. Since these two subspaces are orthogonal and $\mathrm{R}\left(\Delta_{j}\right)$ is closed, $\mathrm{R}\left(P_{j-1}\right)$ and $\mathrm{R}\left(P_{j}^{*}\right)$ are both closed and we have an orthogonal direct sum:

$$
L^{2}\left(V_{j}\right)=\mathrm{N}\left(\Delta_{j}\right) \oplus \mathrm{R}\left(P_{j-1}\right) \oplus \mathrm{R}\left(P_{j}^{*}\right)
$$

This proves (a).
The natural inclusion of $\mathrm{N}\left(\Delta_{j}\right)$ into $\mathrm{N}\left(P_{j}\right)$ defines a natural map of $\mathrm{N}\left(\Delta_{j}\right) \rightarrow H^{j}(P, V)=\mathrm{N}\left(P_{j}\right) / \mathrm{R}\left(P_{j-1}\right)$. Since $\mathrm{R}\left(P_{j-1}\right)$ is orthogonal to $\mathrm{N}\left(\Delta_{j}\right)$, this map is injective. If $f \in C^{\infty}\left(V_{j}\right)$ and $P_{j} f=0$, we can decompose $f=f_{0} \oplus \Delta_{j} f_{1}$ for $f_{0} \in \mathrm{~N}\left(\Delta_{j}\right)$. Since $f$ and $f_{0}$ are smooth, $\Delta_{j} f_{1}$ is smooth so $f_{1} \in C^{\infty}\left(V_{j}\right) . \quad P_{j} f=P_{j} f_{0}+P_{j} \Delta_{j} f_{1}=0$ implies $P_{j} \Delta_{j} f_{1}=0$. We dot this equation with $P_{j} f_{1}$ to conclude:

$$
0=\left(P_{j} f_{1}, P_{j} P_{j}^{*} P_{j} f_{1}+P_{j} P_{j-1} P_{j-1}^{*} f_{1}\right)=\left(P_{j}^{*} P_{j} f_{1}, P_{j}^{*} P_{j} f_{1}\right)
$$

so $P_{j}^{*} P_{j} f_{1}=0$ so $\Delta_{j} f_{1}=P_{j-1} P_{j-1}^{*} f_{1} \in \mathrm{R}\left(P_{j-1}\right)$. This implies $f$ and $f_{0}$ represent the same element of $H^{j}(P, V)$ so the map $\mathrm{N}\left(\Delta_{j}\right) \rightarrow H^{j}(P, V)$ is surjective. This completes the proof.

To illustrate these concepts, we discuss briefly the de Rham complex. Let $T^{*} M$ be the cotangent space of $M$. The exterior algebra $\Lambda\left(T^{*} M\right)$ is the universal algebra generated by $T^{*} M$ subject to the relation $\xi \wedge \xi=0$ for $\xi \in T^{*} M$. If $\left\{e_{1}, \ldots, e_{m}\right\}$ is a basis for $T^{*} M$ and if $I=\left\{1 \leq i_{1}<\right.$ $\left.i_{2}<\cdots<i_{p} \leq m\right\}$, we define $e_{I}=e_{i_{1}} \wedge \cdots \wedge e_{i_{p}}$. The $\left\{e_{I}\right\}$ form a basis for $\Lambda\left(T^{*} M\right)$ which has dimension $2^{m}$. If we define $|I|=p$, then $\Lambda^{p}\left(T^{*} M\right)$ is the span of the $\left\{e_{I}\right\}_{|I|=p}$; this is the bundle of $p$-forms. A section of $C^{\infty}\left(\Lambda^{p} T^{*} M\right)$ is said to be a smooth $p$-form over $M$.

Let $x=\left(x_{1}, \ldots, x_{m}\right)$ be local coordinates on $M$ and let $\left\{d x_{1}, \ldots, d x_{m}\right\}$ be the corresponding frame for $T^{*} M$. If $f \in C^{\infty}(M)=C^{\infty}\left(\Lambda^{0}\left(T^{*} M\right)\right)$, we define:

$$
d f=\sum_{k} \frac{\partial f}{\partial x_{k}} d x_{k}
$$

If $y=\left(y_{1}, \ldots, y_{m}\right)$ is another system of local coordinates on $M$, the identity:

$$
d y_{j}=\sum_{k} \frac{\partial y_{j}}{\partial x_{k}} d x_{k}
$$

means that $d$ is well defined and is independent of the coordinate system chosen. More generally, we define $d\left(f d x_{I}\right)=d f \wedge d x_{I}$ so that, for example,

$$
d\left(\sum_{j} f_{j} d x^{j}\right)=\sum_{j<k}\left\{\frac{\partial f_{k}}{\partial x_{j}}-\frac{\partial f_{j}}{\partial x_{k}}\right\} d x_{j} \wedge d x_{k}
$$

Again this is well defined and independent of the coordinate system. Since mixed partial derivatives commute, $d^{2}=0$ so

$$
d: C^{\infty}\left(\Lambda^{p}\left(T^{*} M\right)\right) \rightarrow C^{\infty}\left(\Lambda^{p+1}\left(T^{*} M\right)\right)
$$

forms a complex.
Let $\xi \in T^{*} M$ and let $\operatorname{ext}(\xi): \Lambda^{p}\left(T^{*} M\right) \rightarrow \Lambda^{p+1}\left(T^{*} M\right)$ be defined by exterior multiplication, i.e., $\operatorname{ext}(\xi) \omega=\xi \wedge \omega$. If we decompose $\xi=\sum \xi_{j} d x_{j}$ relative to a local coordinate frame, then $d f=\sum_{j} \partial f / \partial x_{j} d x_{j}$ implies $\sigma(d)=i \operatorname{ext}(\xi)$; the symbol of exterior differentiation is exterior multiplication up to a factor of $i$. Fix $\xi \neq 0$ and choose a basis $\left\{e_{1}, \ldots, e_{m}\right\}$ for $T^{*} M$ such that $\xi=e_{1}$. Then

$$
\operatorname{ext}(\xi) e_{I}= \begin{cases}0 & \text { if } i_{1}=1 \\ e_{J} \text { for } J=\left\{1, i_{1}, \ldots, i_{p}\right\} & \text { if } i_{1}>1\end{cases}
$$

From this it is clear that $\mathrm{N}(\operatorname{ext}(\xi))=\mathrm{R}(\operatorname{ext}(\xi))$ so the de Rham complex is an elliptic complex.

A Riemannian metric on $M$ defines fiber metrics on $\Lambda^{p}\left(T^{*} M\right)$. If $\left\{e_{i}\right\}$ is an orthonormal local frame for $T M$ we take the dual frame $\left\{e_{i}^{*}\right\}$ for $T^{*} M$. For notational simplicity, we will simply denote this frame again by $\left\{e_{i}\right\}$.

The corresponding $\left\{e_{I}\right\}$ define an orthonormal local frame for $\Lambda\left(T^{*} M\right)$. We define interior multiplication $\operatorname{int}(\xi): \Lambda^{p}\left(T^{*} M\right) \rightarrow \Lambda^{p-1}\left(T^{*} M\right)$ to be the dual of exterior multiplication. Then:

$$
\operatorname{int}\left(e_{1}\right) e_{I}=\left\{\begin{array}{cl}
e_{j} \quad \text { for } J=\left\{i_{2}, \ldots, i_{p}\right\} & \text { id } i_{1}=1 \\
0 & \text { if } i_{1}>1
\end{array}\right.
$$

so

$$
\operatorname{int}\left(e_{1}\right) \operatorname{ext}\left(e_{1}\right)+\operatorname{ext}\left(e_{1}\right) \operatorname{int}\left(e_{1}\right)=I
$$

If $|\xi|^{2}$ denotes the length of the covector $\xi$, then more generally:

$$
(i \operatorname{ext}(\xi)-i \operatorname{int}(\xi))^{2}=|\xi|^{2} I
$$

If $\delta: C^{\infty}\left(\Lambda^{p}\left(T^{*} M\right)\right) \rightarrow C^{\infty}\left(\Lambda^{p-1}\left(T^{*} M\right)\right)$ is the adjoint of $d$, then $\sigma_{L} \delta=$ $i \operatorname{int}(\xi)$. We let $\delta d+d \delta=(d+\delta)^{2}=\Delta . \sigma_{L} \Delta=|\xi|^{2}$ is elliptic. The de Rham theorem gives a natural isomorphism from the cohomology of $M$ to $H^{*}(\Lambda, d)$ :

$$
H^{p}(M ; \mathbf{C})=\mathrm{N}\left(d_{p}\right) / \mathrm{R}\left(d_{p-1}\right)
$$

where we take closed modulo exact forms. The Hodge decomposition theorem implies these groups are naturally isomorphic to the harmonic $p$ forms $\mathrm{N}\left(\Delta_{p}\right)$ which are finite dimensional. The Euler-Poincaré characteristic $\chi(M)$ is given by:

$$
\chi(M)=\sum(-1)^{p} \operatorname{dim} H^{p}(M ; \mathbf{C})=\sum(-1)^{p} \operatorname{dim} \mathrm{~N}\left(\Delta_{p}\right)=\operatorname{index}(d)
$$

is therefore the index of an elliptic complex.
If $M$ is oriented, we let dvol be the oriented volume element. The Hodge * operator $*: \Lambda^{p}\left(T^{*} M\right) \rightarrow \Lambda^{m-p}\left(T^{*} M\right)$ is defined by the identity:

$$
\omega \wedge * \omega=(\omega \cdot \omega) \text { dvol }
$$

where "." denotes the inner product defined by the metric. If $\left\{e_{i}\right\}$ is an oriented local frame for $T^{*} M$, then dvol $=e_{1} \wedge \cdots \wedge e_{m}$ and

$$
*\left(e_{1} \wedge \cdots \wedge e_{p}\right)=e_{p+1} \wedge \cdots \wedge e_{m}
$$

The following identities are immediate consequences of Stoke's theorem:

$$
* *=(-1)^{p(m-p)} \quad \text { and } \quad \delta=(-1)^{m p+m+1} * d * .
$$

Since $\Delta=(d+\delta)^{2}=d \delta+\delta d$ we compute $* \Delta=\Delta *$ so $*: \mathrm{N}\left(\Delta_{p}\right) \rightarrow \mathrm{N}\left(\Delta_{m-p}\right)$ is an isomorphism. We may regard $*$ as an isomorphism $*: H^{p}(M ; \mathbf{C}) \rightarrow$ $H^{m-p}(M ; \mathbf{C})$; in this description it is Poincaré duality.

The exterior algebra is not very suited to computations owing to the large number of signs which enter in the discussion of $*$. When we discuss the signature and spin complexes in the third chapter, we will introduce Clifford algebras which make the discussion of Poincaré duality much easier.

It is possible to "roll up" the de Rham complex and define:

$$
(d+\delta)_{\mathrm{e}}: C^{\infty}\left(\Lambda^{\mathrm{e}}\left(T^{*} M\right)\right) \rightarrow C^{\infty}\left(\Lambda^{\circ}\left(T^{*} M\right)\right)
$$

where

$$
\Lambda^{\mathrm{e}}\left(T^{*} M\right)=\bigoplus_{2 k} \Lambda^{2 k}\left(T^{*} M\right) \quad \text { and } \quad \Lambda^{\circ}\left(T^{*} M\right)=\bigoplus_{2 k+1} \Lambda^{2 k+1}\left(T^{*} M\right)
$$

denote the differential forms of even and odd degrees. $(d+\delta)$ is an elliptic operator since $(d+\delta)_{e}^{*}(d+\delta)_{\mathrm{e}}=\Delta$ is elliptic since $\operatorname{dim} \Lambda^{\mathrm{e}}=\operatorname{dim} \Lambda^{\circ}$. (In this representation $(d+\delta)_{\mathrm{e}}$ is not self-adjoint since the range and domain are distinct, $\left.(d+\delta)_{\mathrm{e}}^{*}=(d+\delta)_{\mathrm{o}}\right)$. It is clear index $(d+\delta)_{\mathrm{e}}=\operatorname{dim} \mathrm{N}\left(\Delta^{\mathrm{e}}\right)=$ $\operatorname{dim} \mathrm{N}\left(\Delta^{\circ}\right)=\chi(M)$. We can always "roll up" any elliptic complex to form an elliptic complex of the same index with two terms. Of course, the original elliptic complex does not depend upon the choice of a fiber metric to define adjoints so there is some advantage in working with the full complex occasionally as we shall see later.

We note finally that if $m$ is even, we can always find a manifold with $\chi(M)$ arbitrary so there exist lots of elliptic operators with non-zero index. We shall see that $\operatorname{index}(P)=0$ if $m$ is odd (and one must consider pseudodifferential operators to get a non-zero index in that case).

We summarize these computations for the de Rham complex in the following:
Lemma 1.5.3. Let $\Lambda(M)=\Lambda\left(T^{*} M\right)$ be the complete exterior algebra.
(a) $d: C^{\infty}\left(\Lambda^{p}\left(T^{*} M\right)\right) \rightarrow C^{\infty}\left(\Lambda^{p+1}\left(T^{*} M\right)\right)$ is an elliptic complex. The symbol is $\sigma_{L}(d)(x, \xi)=i \operatorname{ext}(\xi)$.
(b) If $\delta$ is the adjoint, then $\sigma_{L}(\delta)(x, \xi)=-i \operatorname{int}(\xi)$.
(c) If $\Delta_{p}=(d \delta+\delta d)_{p}$ is the associated Laplacian, then $\mathrm{N}\left(\Delta_{p}\right)$ is finite dimensional and there are natural identifications:

$$
\mathrm{N}\left(\Delta_{p}\right) \simeq \mathrm{N}\left(d_{p}\right) / \mathrm{R}\left(d_{p-1}\right) \simeq H^{p}(M ; \mathbf{C})
$$

(d) $\operatorname{index}(d)=\chi(M)$ is the Euler-Poincaré characteristic of $M$.
(e) If $M$ is oriented, we let * be the Hodge operator. Then

$$
* *=(-1)^{p(m-p)} \quad \text { and } \quad \delta=(-1)^{m p+m+1} * d * .
$$

Furthermore, $*: \mathrm{N}\left(\Delta_{p}\right) \simeq \mathrm{N}\left(\Delta_{m-p}\right)$ gives Poincaré duality.
In the next chapter, we will discuss the Gauss-Bonnet theorem which gives a formula for index $(d)=\chi(M)$ in terms of curvature.

### 1.6. The Heat Equation.

Before proceeding with our discussion of the index of an elliptic operator, we must discuss spectral theory. We restrict ourselves to the context of a compact self-adjoint operator to avoid unnecessary technical details. Let $T \in \operatorname{COM}(H)$ be a self-adjoint compact operator on the Hilbert space $H$. Let

$$
\operatorname{spec}(T)=\{\lambda \in \mathbf{C}:(T-\lambda) \notin \operatorname{GL}(H)\}
$$

Since $\mathrm{GL}(H)$ is open, $\operatorname{spec}(T)$ is a closed subset of $C$. If $|\lambda|>|T|$, the series

$$
g(\lambda)=\sum_{n=0}^{\infty} T^{n} / \lambda^{n+1}
$$

converges to define an element of $\operatorname{END}(H)$. As $(T-\lambda) g(\lambda)=g(\lambda)(T-\lambda)=$ $-I, \lambda \notin \operatorname{spec}(T)$. This shows $\operatorname{spec}(T)$ is bounded.

Since $T$ is self-adjoint, $\mathrm{N}(T-\bar{\lambda})=\{0\}$ if $\lambda \notin \mathbf{R}$. This implies $\mathrm{R}(T-\lambda)$ is dense in $H$. Since $T x \cdot x$ is always real, $|(T-\lambda) x \cdot x| \geq\left.\operatorname{im}(\lambda)| | x\right|^{2}$ so that $|(T-\lambda) x| \geq \operatorname{im}(\lambda)| | x \mid$. This implies $\mathrm{R}(T-\lambda)$ is closed so $T-\lambda$ is surjective. Since $T-\lambda$ is injective, $\lambda \notin \operatorname{spec}(T)$ if $\lambda \in \mathbf{C}-\mathbf{R} \operatorname{so} \operatorname{spec}(T)$ is a subset of $\mathbf{R}$. If $\lambda \in[-|T|,|T|]$, we define $E(\lambda)=\{x \in H: T x=\lambda x\}$.

Lemma 1.6.1. Let $T \in \operatorname{COM}(H)$ be self-adjoint. Then

$$
\operatorname{dim}\{E(-|T|) \oplus E(|T|)\}>0 .
$$

Proof: If $|T|=0$ then $T=0$ and the result is clear. Otherwise choose $\left|x_{n}\right|=1$ so that $\left|T x_{n}\right| \rightarrow|T|$. We choose a subsequence so $T x_{n} \rightarrow y$. Let $\lambda=|T|$, we compute:

$$
\begin{aligned}
\left|T^{2} x_{n}-\lambda^{2} x_{n}\right|^{2} & =\left|T^{2} x_{n}\right|^{2}+\left|\lambda^{2} x_{n}\right|^{2}-2 \lambda^{2} T^{2} x_{n} \cdot x_{n} \\
& \leq 2 \lambda^{4}-2 \lambda^{2}\left|T x_{n}\right|^{2} \rightarrow 0 .
\end{aligned}
$$

Since $T x_{n} \rightarrow y, \quad T^{2} x_{n} \rightarrow T y$. Thus $\lambda^{2} x_{n} \rightarrow T y$. Since $\lambda^{2} \neq 0$, this implies $x_{n} \rightarrow x$ for $x=T y / \lambda^{2}$. Furthermore $\left|T^{2} x-\lambda^{2} x\right|=0$. Since $\left(T^{2}-\lambda^{2}\right)=(T-\lambda)(T+\lambda)$, either $(T+\lambda) x=0$ so $x \in E(-\lambda) \neq\{0\}$ or $(T-\lambda)\left(y_{1}\right)=0$ for $y_{1}=(T+\lambda) x \neq 0$ so $E(\lambda) \neq\{0\}$. This completes the proof.

If $\lambda \neq 0$, the equation $T x=\lambda x$ implies the unit disk in $E(\lambda)$ is compact and hence $E(\lambda)$ is finite dimensional. $E(\lambda)$ is $T$-invariant. Since $T$ is selfadjoint, the orthogonal complement $E(\lambda)$ is also $T$-invariant. We take an orthogonal decomposition $H=E(|T|) \oplus E(-|T|) \oplus H_{1}$. $\quad T$ respects this decomposition; we let $T_{1}$ be the restriction of $T$ to $H_{1}$. Clearly $\left|T_{1}\right| \leq|T|$. If we have equality, then Lemma 1.6.1 implies there exists $x \neq 0$ in $H_{1}$
so $T_{1} x= \pm|T| x$, which would be false. Thus $\left|T_{1}\right|<|T|$. We proceed inductively to decompose:

$$
H=E\left(\lambda_{1}\right) \oplus E\left(-\lambda_{1}\right) \oplus \cdots \oplus E\left(\lambda_{n}\right) \oplus E\left(-\lambda_{n}\right) \oplus H_{n}
$$

where $\lambda_{1}>\lambda_{2}>\cdots>\lambda_{n}$ and where $\left|T_{n}\right|<\lambda_{n}$. We also assume $E\left(\lambda_{n}\right) \oplus$ $E\left(-\lambda_{n}\right) \neq\{0\}$. We suppose that the $\lambda_{n}$ do not converge to zero but converge to $\varepsilon$ positive. For each $n$, we choose $x_{n}$ so $\left|x_{n}\right|=1$ and $T x_{n}=$ $\pm \lambda_{n} x_{n}$. Since $T$ is self-adjoint, this is an orthogonal decomposition; $\mid x_{j}-$ $x_{k} \mid=\sqrt{2}$. Since $T$ is compact, we can choose a convergent subsequence $\lambda_{n} x_{n} \rightarrow y$. Since $\lambda_{n} \rightarrow \varepsilon$ positive, this implies $x_{n} \rightarrow x$. This is impossible so therefore the $\lambda_{n} \rightarrow 0$. We define $H_{0}=\bigcap_{n} H_{n}$ as a closed subset of $H$. Since $|T|<\lambda_{n}$ for all $n$ on $H_{0},|T|=0$ so $T=0$ and $H_{0}=E(0)$. This defines a direct sum decomposition of the form:

$$
H=\bigoplus_{k} E\left(\mu_{k}\right) \oplus E(0)
$$

where the $\mu_{k} \in \mathbf{R}$ are the non-zero subspace of the $E\left(\lambda_{n}\right)$ and $E\left(-\lambda_{n}\right)$. We construct a complete orthonormal system $\left\{\phi_{n}\right\}_{n=1}^{\infty}$ for $H$ with either $\phi_{n} \in E(0)$ or $\phi_{n} \in E\left(\mu_{k}\right)$ for some $k$. Then $T \phi_{n}=\lambda_{n} \phi_{n}$ so $\phi_{n}$ is an eigenvector of $T$. This proves the spectral decomposition for self-adjoint compact operators:

Lemma 1.6.2. Let $T \in \operatorname{COM}(H)$ be self-adjoint. We can find a complete orthonormal system for $H$ consisting of eigenvectors of $T$.

We remark that this need not be true if $T$ is self-adjoint but not compact or if $T$ is compact but not self-adjoint.

We can use this lemma to prove the following:
Lemma 1.6.3. Let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be an elliptic self-adjoint $\Psi D O$ of order $d>0$.
(a) We can find a complete orthonormal basis $\left\{\phi_{n}\right\}_{n=1}^{\infty}$ for $L^{2}(V)$ of eigenvectors of $P . \quad P \phi_{n}=\lambda_{n} \phi_{n}$.
(b) The eigenvectors $\phi_{n}$ are smooth and $\lim _{n \rightarrow \infty}\left|\lambda_{n}\right|=\infty$.
(c) If we order the eigenvalues $\left|\lambda_{1}\right| \leq\left|\lambda_{2}\right| \leq \cdots$ then there exists a constant $C>0$ and an exponent $\delta>0$ such that $\left|\lambda_{n}\right| \geq C n^{\delta}$ if $n>n_{0}$ is large.
Remark: The estimate (c) can be improved to show $\left|\lambda_{n}\right| \sim n^{d / m}$ but the weaker estimate will suffice and is easier to prove.

Proof: $P: H_{d}(V) \rightarrow L^{2}(V)$ is Fredholm. $P: \mathrm{N}(P)^{\perp} \cap H_{d}(M) \rightarrow \mathrm{N}(P)^{\perp} \cap$ $L^{2}(V)$ is 1-1 and onto; by Gärding's inequality $P \phi \in L^{2}$ implies $\phi \in H_{d}$. We let $S$ denote the inverse of this map and extend $S$ to be zero on the finite dimensional space $\mathrm{N}(P)$. Since the inclusion of $H_{d}(M)$ into $L^{2}(V)$ is compact, $S$ is a compact self-adjoint operator. $S$ is often referred to
as the Greens operator. We find a complete orthonormal system $\left\{\phi_{n}\right\}$ of eigenvectors of $S$. If $S \phi_{n}=0$ then $P \phi_{n}=0$ since $\mathrm{N}(S)=\mathrm{N}(P)$. If $S \phi_{n}=\mu_{n} \phi_{n} \neq 0$ then $P \phi_{n}=\lambda_{n} \phi_{n}$ for $\lambda_{n}=\mu_{n}^{-1}$. Since the $\mu_{n} \rightarrow 0$, the $\left|\lambda_{n}\right| \rightarrow \infty$. If $k$ is an integer such that $d k>1$, then $P^{k}-\lambda_{n}^{k}$ is elliptic. Since $\left(P^{k}-\lambda_{n}^{k}\right) \phi_{n}=0$, this implies $\phi_{n} \in C^{\infty}(V)$. This completes the proof of (a) and (b).

By replacing $P$ with $P^{k}$ we replace $\lambda_{n}$ by $\lambda_{n}^{k}$. Since $k d>\frac{m}{2}$ if $k$ is large, we may assume without loss of generality that $d>\frac{m}{2}$ in the proof of (c). We define:

$$
|f|_{\infty, 0}=\sup _{x \in M}|f(x)| \quad \text { for } f \in C^{\infty}(V)
$$

We estimate:

$$
|f|_{\infty, 0} \leq C|f|_{d} \leq C\left(|P f|_{0}+|f|_{0}\right)
$$

Let $F(a)$ be the space spanned by the $\phi_{j}$ where $\left|\lambda_{j}\right| \leq a$ and let $n=$ $\operatorname{dim} F(a)$. We estimate $n=n(a)$ as follows. We have

$$
|f|_{\infty, 0} \leq C(1+a)|f|_{0} \quad \text { on } F(a)
$$

Suppose first $V=M \times \mathbf{C}$ is the trivial line bundle. Let $c_{j}$ be complex constants, then this estimate shows:

$$
\left|\sum_{j=1}^{n} c_{j} \phi_{j}(x)\right| \leq C(1+a)\left\{\sum_{j=1}^{n}\left|c_{j}\right|^{2}\right\}^{1 / 2}
$$

If we take $c_{j}=\bar{\phi}_{j}(x)$ then this yields the estimate:

$$
\sum_{j=1}^{n} \phi_{j}(x) \bar{\phi}_{j}(x) \leq C(1+a)\left\{\sum_{j=1}^{n} \phi_{j}(x) \bar{\phi}_{j}(x)\right\}^{1 / 2}
$$

i.e.,

$$
\sum_{j=1}^{n} \phi_{j}(x) \bar{\phi}_{j}(x) \leq C^{2}(1+a)^{2}
$$

We integrate this estimate over $M$ to conclude

$$
n \leq C^{2}(1+a)^{2} \operatorname{vol}(M)
$$

or equivalently

$$
C_{1}\left(n-C_{2}\right)^{1 / 2} \leq a=\left|\lambda_{n}\right|
$$

from which the desired estimate follows.

If $\operatorname{dim} V=k$, we choose a local orthonormal frame for $V$ to decompose $\phi_{j}$ into components $\phi_{j}^{u}$ for $1 \leq u \leq k$. We estimate:

$$
\left|\sum_{j=1}^{n} c_{j}^{u} \phi_{j}^{u}(x)\right| \leq C(1+a)\left\{\sum_{j=1}^{n}\left|c_{j}^{u}\right|^{2}\right\}^{1 / 2} \quad \text { for } u=1, \ldots, k
$$

If we let $c_{j}^{u}=\bar{\phi}_{j}^{u}(x)$, then summing over $u$ yields:

$$
\left|\sum_{j=1}^{n} \phi_{j} \bar{\phi}_{j}(x)\right| \leq k C(1+a)\left|\sum_{j=1}^{n} \phi_{j} \bar{\phi}_{j}(x)\right|^{1 / 2}
$$

since each term on the left hand side of the previous inequality can be estimated seperately by the right hand side of this inequality. This means that the constants which arise in the estimate of (c) depend on $k$, but this causes no difficulty. This completes the proof of (c).

The argument given above for (c) was shown to us by Prof. B. Allard (Duke University) and it is a clever argument to avoid the use of the Schwarz kernel theorem.

Let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be an elliptic $\Psi D O$ of order $d>0$ which is self-adjoint. We say $P$ has positive definite leading symbol if there exists $p(x, \xi): T^{*} M \rightarrow \operatorname{END}(V)$ such that $p(x, \xi)$ is a positive definite Hermitian matrix for $\xi \neq 0$ and such that $\sigma P-p \in S^{d-1}$ in any coordinate system. The spectrum of such a $P$ is not necessarily non-negative, but it is bounded from below as we shall show. We construct $Q_{0}$ with leading symbol $\sqrt{p}$ and let $Q=Q_{0}^{*} Q_{0}$. Then by hypothesis $P-Q \in S^{d-1}$. We compute:

$$
\begin{aligned}
(P f, f) & =((P-Q) f, f)+(Q f, f) \\
|((P-Q) f, f)| & \leq C|f|_{d / 2}|(P-Q) f|_{d / 2} \leq C|f|_{d / 2}|f|_{d / 2-1} \\
(Q f, f) & =\left(Q_{0} f, Q_{0} f\right) \quad \text { and } \quad|f|_{d / 2}^{2}<C|f|_{0}^{2}+C\left|Q_{0} f\right|_{0}^{2}
\end{aligned}
$$

We use this to estimate for any $\varepsilon>0$ that:

$$
\begin{aligned}
|((P-Q) f, f)| & \leq C|f|_{d / 2}\left(|f|_{d / 2-1}\right) \leq \varepsilon|f|_{d / 2}^{2}+C(\varepsilon)|f|_{d / 2}|f|_{0} \\
& \leq 2 \varepsilon|f|_{d / 2}^{2}+C(\varepsilon)|f|_{0}^{2} \\
& \leq 2 C \varepsilon\left|Q_{0} f\right|_{0}^{2}+C(\varepsilon)|f|_{0}^{2} .
\end{aligned}
$$

Choose $\varepsilon$ so $2 C \varepsilon \leq 1$ and estimate:

$$
(P f, f) \geq\left(Q_{0} f, Q_{0} f\right)-|((P-Q) f, f)| \geq-C(\varepsilon)|f|_{0}^{2}
$$

This implies:

Lemma 1.6.4. Let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be an elliptic $\Psi D O$ of order $d>0$ which is self-adjoint with positive definite leading symbol. Then $\operatorname{spec}(P)$ is contained in $[-C, \infty)$ for some constant $C$.

We fix such a $P$ henceforth. The heat equation is the partial differential equation:

$$
\left(\frac{d}{d t}+P\right) f(x, t)=0 \quad \text { for } t \geq 0 \text { with } f(x, 0)=f(x)
$$

At least formally, this has the solution $f(x, t)=e^{-t P} f(x)$. We decompose $f(x)=\sum c_{n} \phi_{n}(x)$ for $c_{n}=\left(f, \phi_{n}\right)$ in a generalized Fourier series. The solution of the heat equation is given by:

$$
f(x, t)=\sum_{n} e^{-t \lambda_{n}} c_{n} \phi_{n}(x)
$$

Proceeding formally, we define:

$$
K(t, x, y)=\sum_{n} e^{-t \lambda_{n}} \phi_{n}(x) \otimes \bar{\phi}_{n}(y): V_{y} \rightarrow V_{x}
$$

so that:

$$
\begin{aligned}
e^{-t P} f(x) & =\int_{M} K(t, x, y) f(y) \operatorname{dvol}(y) \\
& =\sum_{n} e^{-t \lambda_{n}} \phi_{n}(x) \int_{M} f(y) \cdot \phi_{n}(y) \operatorname{dvol}(y)
\end{aligned}
$$

We regard $K(t, x, y)$ as an endomorphism from the fiber of $V$ over $y$ to the fiber of $V$ over $x$.

We justify this purely formal procedure using Lemma 1.3.4. We estimate that:

$$
\left|\phi_{n}\right|_{\infty, k} \leq C\left(\left|\phi_{n}\right|_{0}+\left|P^{j} \phi\right|_{0}\right)=C\left(1+\left|\lambda_{n}\right|^{j}\right) \quad \text { where } j d>k+\frac{m}{2}
$$

Only a finite number of eigenvalues of $P$ are negative by Lemma 1.6.4. These will not affect convergence questions, so we may assume $\lambda>0$. Estimate:

$$
e^{-t \lambda} \lambda^{j} \leq t^{-j} C(j) e^{-t \lambda / 2}
$$

to compute:

$$
|K(t, x, y)|_{\infty, k} \leq t^{-j(k)} C(k) \sum_{n} e^{-t \lambda_{n} / 2}
$$

Since $|\lambda| \geq C n^{\delta}$ for $\delta>0$ and $n$ large, the series can be bounded by

$$
\sum_{n>0} e^{-t n^{\delta / 2}}
$$

which converges. This shows $K(t, x, y)$ is an infinitely smooth function of $(t, x, y)$ for $t>0$ and justifies all the formal procedures involved. We compute

$$
\operatorname{Tr}_{L^{2}} e^{-t P}=\sum e^{-t \lambda_{n}}=\int_{M} \operatorname{Tr}_{V_{x}} K(t, x, x) \operatorname{dvol}(x)
$$

We can use this formula to compute the index of $Q$ :
Lemma 1.6.5. Let $Q: C^{\infty}(V) \rightarrow C^{\infty}(W)$ be an elliptic $\Psi D O$ of order $d>0$. Then for $t>0, e^{-t Q^{*} Q}$ and $e^{-t Q Q^{*}}$ are in $\Psi_{-\infty}$ with smooth kernel functions and

$$
\operatorname{index}(Q)=\operatorname{Tr} e^{-t Q^{*} Q}-\operatorname{Tr} e^{-t Q Q^{*}} \quad \text { for any } t>0
$$

Proof: Since $e^{-t Q^{*} Q}$ and $e^{-t Q Q^{*}}$ have smooth kernel functions, they are in $\Psi_{-\infty}$ so we must only prove the identity on index $(Q)$. We define $E_{0}(\lambda)=$ $\left\{\phi \in L^{2}(V): Q^{*} Q \phi=\lambda \phi\right\}$ and $E_{1}(\lambda)=\left\{\phi \in L^{2}(W): Q Q^{*} \phi=\lambda \phi\right\}$. These are finite dimensional subspaces of smooth sections to $V$ and $W$. Because $Q\left(Q^{*} Q\right)=\left(Q Q^{*}\right) Q$ and $Q^{*}\left(Q Q^{*}\right)=\left(Q^{*} Q\right) Q^{*}, Q$ and $Q^{*}$ define maps:

$$
Q: E_{0}(\lambda) \rightarrow E_{1}(\lambda) \quad \text { and } \quad Q^{*}: E_{1}(\lambda) \rightarrow E_{0}(\lambda)
$$

If $\lambda \neq 0, \lambda=Q^{*} Q: E_{0}(\lambda) \rightarrow E_{1}(\lambda) \rightarrow E_{0}(\lambda)$ is an isomorphism so $\operatorname{dim} E_{0}(\lambda)=\operatorname{dim} E_{1}(\lambda)$. We compute:

$$
\begin{aligned}
\operatorname{Tr} e^{-t Q^{*} Q}-\operatorname{Tr} e^{-t Q Q^{*}} & =\sum_{\lambda} e^{-t \lambda}\left\{\operatorname{dim} E_{0}(\lambda)-\operatorname{dim} E_{1}(\lambda)\right\} \\
& =e^{-t 0}\left\{\operatorname{dim} E_{0}(0)-\operatorname{dim} E_{1}(0)\right\} \\
& =\operatorname{index}(Q)
\end{aligned}
$$

which completes the proof.
If $(Q, V)$ is a elliptic complex, we use the same reasoning to conclude $e^{-t \Delta_{i}}$ is in $\Psi_{-\infty}$ with a smooth kernel function. We define $E_{i}(\lambda)=\{\phi \in$ $\left.L^{2}\left(V_{i}\right): \Delta_{i} \phi=\lambda \phi\right\}$. Then $Q_{i}: E_{i}(\lambda) \rightarrow E_{i+1}(\lambda)$ defines an acyclic complex (i.e., $\mathrm{N}\left(Q_{i}\right)=\mathrm{R}\left(Q_{i-1}\right)$ ) if $\lambda \neq 0$ so that $\sum(-1)^{i} \operatorname{dim} E_{i}(\lambda)=0$ for $\lambda \neq 0$. Therefore

$$
\operatorname{index}(Q)=\sum_{i}(-1)^{i} \operatorname{Tr}\left(e^{-t \Delta_{i}}\right)
$$

and Lemma 1.6.5 generalizes to the case of elliptic complexes which are not just two term complexes.

Let $P$ be an elliptic $\Psi$ DO of order $d>0$ which is self-adjoint with positive definite leading symbol. Then $\operatorname{spec}(P)$ is contained in $[-C, \infty)$ for some constant $C$. For $\lambda \notin \operatorname{spec}(P),(P-\lambda)^{-1} \in \operatorname{END}\left(L^{2}(V)\right)$ satisfies:

$$
\left|(P-\lambda)^{-1}\right|=\operatorname{dist}(\lambda, \operatorname{spec}(P))^{-1}=\inf _{n}\left|\lambda-\lambda_{n}\right|^{-1}
$$

The function $\left|(P-\lambda)^{-1}\right|$ is a continuous function of $\lambda \in \mathbf{C}-\operatorname{spec}(P)$. We use the Riemann integral with values in the Banach space $\operatorname{END}\left(L^{2}(V)\right)$ to define:

$$
e^{-t P}=\frac{1}{2 \pi i} \int_{\gamma} e^{-t \lambda}(P-\lambda)^{-1} d \lambda
$$

where $\gamma$ is the path about $[-C, \infty)$ given by the union of the appropriate pieces of the straight lines $\operatorname{Re}(\lambda+C+1)= \pm \operatorname{Im}(\lambda)$ pictured below. We let $\mathcal{R}$ be the closed region of $\mathbf{C}$ consisting of $\gamma$ together with that component of $\mathbf{C}-\gamma$ which does not contain $[-C, \infty)$.


We wish to extend $(P-\lambda)^{-1}$ to $H_{s}$. We note that

$$
|\lambda-\mu|^{-1} \leq C \quad \text { for } \lambda \in \mathcal{R}, \mu \in \operatorname{spec}(P)
$$

so that $\left|(P-\lambda)^{-1} f\right|_{0}<C|f|_{0}$. We use Lemma 1.3.5 to estimate:

$$
\begin{aligned}
\left|(P-\lambda)^{-1} f\right|_{k d} & \leq C\left\{\left|P^{k}(P-\lambda)^{-1} f\right|_{0}+\left|(P-\lambda)^{-1} f\right|_{0}\right\} \\
& \leq C\left\{\left|P^{k-1} f\right|_{0}+\left|\lambda P^{k-1}(P-\lambda)^{-1} f\right|_{0}+|f|_{0}\right\} \\
& \leq C\left\{|f|_{k d-d}+|\lambda|\left|(P-\lambda)^{-1} f\right|_{k d-d}\right\}
\end{aligned}
$$

If $k=1$, this implies $\left|(P-\lambda)^{-1} f\right|_{d} \leq C(1+|\lambda|)|f|_{0}$. We now argue by induction to estimate:

$$
\left|(P-\lambda)^{-1} f\right|_{k d} \leq C(1+|\lambda|)^{k-1}|f|_{k d-d}
$$

We now interpolate. If $s>0$, choose $k$ so $k d \geq s>k d-d$ and estimate:

$$
\begin{aligned}
\left|(P-\lambda)^{-1} f\right|_{s} & \leq C\left|(P-\lambda)^{-1} f\right|_{k d} \leq C(1+|\lambda|)^{k-1}|f|_{k d-d} \\
& \leq C(1+|\lambda|)^{k-1}|f|_{s}
\end{aligned}
$$

Similarly, if $s<0$, we use duality to estimate:

$$
\left|\left((P-\lambda)^{-1} f, g\right)\right|=\left|\left(f,(P-\bar{\lambda})^{-1} g\right)\right| \leq C(1+|\lambda|)^{k-1}|f|_{s}|g|_{-s}
$$

which by Lemma 1.3.4 shows

$$
\left|(P-\lambda)^{-1} f\right|_{s}<C(1+|\lambda|)^{k-1}|f|_{s}
$$

in this case as well.
Lemma 1.6.6. Let $P$ be a elliptic $\Psi D O$ of order $d>0$ which is selfadjoint. We suppose the leading symbol of $P$ is positive definite. Then:
(a) Given $s$ there exists $k=k(s)$ and $C=C(s)$ so that

$$
\left|(P-\lambda)^{-1} f\right|_{s} \leq C(1+|\lambda|)^{k}|f|_{s}
$$

for all $\lambda \in \mathcal{R}$.
(b) Given $j$ there exists $k=k(j, d)$ so $\left(P^{k}-\lambda\right)^{-1}$ represents a smoothing operator with $C^{j}$ kernel function which is of trace class for any $\lambda \in \mathcal{R}$.

Proof: (a) follows from the estimates previously. To prove (b) we let

$$
K_{k}(x, y, \lambda)=\sum_{n} \frac{1}{\lambda_{n}^{k}-\lambda} \phi_{n}(x) \otimes \phi_{n}(y)
$$

be the kernel of $\left(P^{k}-\lambda\right)^{-1}$. The region $\mathcal{R}$ was chosen so $|\tilde{\lambda}-\lambda| \geq \varepsilon|\tilde{\lambda}|$ for some $\varepsilon>0$ and $\tilde{\lambda} \in \mathbf{R}^{+}$. Therefore $\left|\left(\lambda_{n}^{k}-\lambda\right)\right|^{-1} \leq \varepsilon\left|\lambda_{n}^{k}\right|^{-1} \leq \varepsilon^{-1} n^{-k \delta}$ by Lemma 1.6.3. The convergence of the sum defining $K_{k}$ then follows using the same arguments as given in the proof that $e^{-t P}$ is a smoothing operator.

This technical lemma will be used in the next subsection to estimate various error terms which occur in the construction of a parametrix.

### 1.7. Local Formula for <br> The Index of an Elliptic Operator.

In this subsection, let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be a self-adjoint elliptic partial differential operator of order $d>0$. Decompose the symbol $\sigma(P)=$ $p_{d}+\cdots+p_{0}$ into homogeneous polymonials $p_{j}$ of order $j$ in $\xi \in T^{*} M$. We assume $p_{d}$ is a positive definite Hermitian matrix for $\xi \neq 0$. Let the curve $\gamma$ and the region $\mathcal{R}$ be as defined previously.

The operator $(P-\lambda)^{-1}$ for $\lambda \in \mathcal{R}$ is not a pseudo-differential operator. We will approximate $(P-\lambda)^{-1}$ by a pseudo-differential operator $R(\lambda)$ and then use that approximation to obtain properties of $\exp (-t P)$. Let $U$ be an open subset of $\mathbf{R}^{m}$ with compact closure. Fix $d \in \mathbf{Z}$. We make the following definition to generalize that given in section 1.2:
Definition. $q(x, \xi, \lambda) \in S^{k}(\lambda)(U)$ is a symbol of order $k$ depending on the complex parameter $\lambda \in \mathcal{R}$ if
(a) $q(x, \xi, \lambda)$ is smooth in $(x, \xi, \lambda) \in \mathbf{R}^{m} \times \mathbf{R}^{m} \times \mathcal{R}$, has compact $x$-support in $U$ and is holomorphic in $\lambda$.
(b) For all $(\alpha, \beta, \gamma)$ there exist constants $C_{\alpha, \beta, \gamma}$ such that:

$$
\left|D_{x}^{\alpha} D_{\xi}^{\beta} D_{\lambda}^{\gamma} q(x, \xi, \gamma)\right| \leq C_{\alpha, \beta, \gamma}\left(1+|\xi|+|\lambda|^{1 / d}\right)^{k-|\beta|-d|\gamma|} .
$$

We say that $q(x, \xi, \gamma)$ is homogeneous of order $k$ in $(\xi, \lambda)$ if

$$
q\left(x, t \xi, t^{d} \lambda\right)=t^{k} q(x, \xi, \gamma) \quad \text { for } t \geq 1
$$

We think of the parameter $\lambda$ as being of order $d$. It is clear if $q$ is homogeneous in $(\xi, \lambda)$, then it satisfies the decay conditions (b). Since the $p_{j}$ are polynomials in $\xi$, they are regular at $\xi=0$ and define elements of $S^{j}(\lambda)$. If the $p_{j}$ were only pseudo-differential, we would have to smooth out the singularity at $\xi=0$ which would destroy the homogeneity in $(\xi, \lambda)$ and they would not belong to $S^{j}(\lambda)$; equivalently, $d_{\xi}^{\alpha} p_{j}$ will not exhibit decay in $\lambda$ for $|\alpha|>j$ if $p_{j}$ is not a polynomial. Thus the restriction to differential operators is an essential one, although it is possible to discuss some results in the pseudo-differential case by taking more care with the estimates involved.

We also note that $\left(p_{d}-\lambda\right)^{-1} \in S^{-d}(\lambda)$ and that the spaces $S^{*}(\lambda)$ form a symbol class closed under differentiation and multiplication. They are suitable generalizations of ordinary pseudo-differential operators discussed earlier.

We let $\Psi_{k}(\lambda)(U)$ be the set of all operators $Q(\lambda): C_{0}^{\infty}(U) \rightarrow C_{0}^{\infty}(U)$ with symbols $q(x, \xi, \gamma)$ in $S^{k}(\lambda)$ having $x$-support in $U$; we let $q=\sigma Q$. For any fixed $\lambda, Q(\lambda) \in \Psi_{k}(U)$ is an ordinary pseudo-differential operator of order $k$. The new features arise from the dependence on the parameter $\lambda$. We extend the definition of $\sim$ given earlier to this wider class by:

$$
q \sim \sum_{j} q_{j}
$$

if for every $k>0$ there exists $n(k)$ so $n \geq n(k)$ implies $q-\sum_{j \leq n} q_{j} \in$ $S^{-k}(\lambda)$. Lemmas 1.2.1, 1.2.3(b), and 1.3 .2 generalize easily to yield the following Lemma. We omit the proofs in the interests of brevity as they are identical to those previously given with only minor technical modifications.

Lemma 1.7.1.
(a) Let $Q_{i} \in \Psi_{k_{i}}(\lambda)(U)$ with symbol $q_{i}$. Then $Q_{1} Q_{2} \in \Psi_{k_{1}+k_{2}}(\lambda)(U)$ has symbol $q$ where:

$$
q \sim \sum_{\alpha} d_{\xi}^{\alpha} q_{1} D_{x}^{\alpha} q_{2} / \alpha!
$$

(b) Given $n>0$ there exists $k(n)>0$ such that $Q \in \Psi_{-k(n)}(\lambda)(U)$ implies $Q(\lambda)$ defines a continuous map from $H_{-n} \rightarrow H_{n}$ and we can estimate the operator norm by:

$$
|Q(\lambda)|_{-n, n} \leq C(1+|\lambda|)^{-n}
$$

(c) If $h: U \rightarrow \widetilde{U}$ is a diffeomorphism, then $h_{*}: \Psi_{k}(\lambda)(U) \rightarrow \Psi_{k}(\lambda)(\widetilde{U})$ and

$$
\sigma\left(h_{*} P\right)-p\left(h^{-1} x_{1},\left(d h^{-1}\left(x_{1}\right)\right)^{t} \xi_{1}, \lambda\right) \in S^{k-1}(\lambda)(\widetilde{U})
$$

As before, we let $\Psi(\lambda)(U)=\bigcup_{n} \Psi_{n}(\lambda)(U)$ be the set of all pseudodifferential operators depending on a complex parameter $\lambda \in \mathcal{R}$ defined over $U$. This class depends on the order $d$ chosen and also on the region $\mathcal{R}$, but we supress this dependence in the interests of notational clarity. There is no analogue of the completeness of Lemma 1.2.8 for such symbols since we require analyticity in $\lambda$. Thus in constructing an approximation to the parametrix, we will always restrict to a finite sum rather than an infinite sum.

Using (c), we extend the class $\Psi(\lambda)$ to compact manifolds using a partition of unity argument. (a) and (b) generalize suitably. We now turn to the question of ellipticity. We wish to solve the equation:

$$
\sigma(R(\lambda)(P-\lambda))-I \sim 0
$$

Inductively we define $R(\lambda)$ with symbol $r_{0}+r_{1}+\cdots$ where $r_{j} \in S^{-d-j}(\lambda)$. We define $p_{j}^{\prime}(x, \xi, \lambda)=p_{j}(x, \xi)$ for $j<d$ and $p_{d}^{\prime}(x, \xi, \lambda)=p_{d}(x, \xi)-\lambda$. Then $\sigma(P-\lambda)=\sum_{j=0}^{d} p_{j}^{\prime}$. We note that $p_{j}^{\prime} \in S^{j}(\lambda)$ and that $p_{d}^{\prime-1} \in$ $S^{-d}(\lambda)$ so that $(P-\lambda)$ is elliptic in a suitable sense. The essential feature of this construction is that the parameter $\lambda$ is absorbed in the leading symbol and not treated as a lower order perturbation.

The equation $\sigma(R(\lambda)(P-\lambda)) \sim I$ yields:

$$
\sum_{\alpha, j, k} d_{\xi}^{\alpha} r_{j} \cdot D_{x}^{\alpha} p_{k}^{\prime} / \alpha!\sim I
$$

We decompose this series into terms homogeneous of order $-n$ to write:

$$
\sum_{n} \sum_{|\alpha|+j+d-k=n} d_{\xi}^{\alpha} r_{j} \cdot D_{x}^{\alpha} p_{k}^{\prime} / \alpha!\sim I
$$

where $j, k \geq 0$ and $k \leq d$. There are no terms with $n<0$, i.e., positive order. If we investigate the term with $n=0$, we arrive at the condition $r_{0} p_{d}^{\prime}=I$ so $r_{0}=\left(p_{d}-\lambda\right)^{-1}$ and inductively:

$$
r_{n}=-r_{0} \sum_{\substack{|\alpha|+j+d-k=n \\ j<n}} d_{\xi}^{\alpha} r_{j} D_{x}^{\alpha} p_{k}^{\prime} / \alpha!
$$

If $k=d$ then $|\alpha|>0$ so $D_{x}^{\alpha} p_{k}^{\prime}=D_{x}^{\alpha} p_{k}$ in this sum. Therefore we may replace $p_{k}^{\prime}$ by $p_{k}$ and write

$$
r_{n}=-r_{0} \sum_{\substack{|\alpha|+j+d-k=n \\ j<n}} d_{\xi}^{\alpha} r_{j} D_{x}^{\alpha} p_{k} / \alpha!
$$

In a similar fashion, we can define $\widetilde{R}(\lambda)$ so $\sigma((P-\lambda) \widetilde{R}(\lambda)-I) \sim 0$. This implies $\sigma(R(\lambda)-\widetilde{R}(\lambda)) \sim 0$ so $R(\lambda)$ provides a formal left and right inverse.

Since such an inverse is unique modulo lower order terms, $R(\lambda)$ is well defined and unique modulo lower order terms in any coordinate system. We define $R(\lambda)$ globally on $M$ using a partition of unity argument. To avoid convergence questions, we shall let $R(\lambda)$ have symbol $r_{0}+\cdots+r_{n_{0}}$ where $n_{0}$ is chosen to be very large. $R(\lambda)$ is unique modulo $\Psi_{-n_{0}-d}(\lambda)$. For notational convenience, we supress the dependence of $R(\lambda)$ upon $n_{0}$.
$R(\lambda)$ gives a good approximation to $(P-\lambda)^{-1}$ in the following sense:
Lemma 1.7.2. Let $k>0$ be given. We can choose $n_{0}=n_{0}(k)$ so that
$\left|\left\{(P-\lambda)^{-1}-R(\lambda)\right\} f\right|_{k} \leq C_{k}(1+|\lambda|)^{-k}|f|_{-k} \quad$ for $\lambda \in \mathcal{R}, \quad f \in C^{\infty}(V)$.
Thus $(P-\lambda)^{-1}$ is approximated arbitrarily well by the parametrix $R(\lambda)$ in the operator norms as $\lambda \rightarrow \infty$.

Proof: We compute:

$$
\begin{aligned}
\left|\left\{(P-\lambda)^{-1}-R(\lambda)\right\} f\right|_{k} & =\left|(P-\lambda)^{-1}\{I-(P-\lambda) R(\lambda)\} f\right|_{k} \\
& \leq C_{k}\left(1+|\lambda|^{\nu}\right)|(I-(P-\lambda) R(\lambda)) f|_{k}
\end{aligned}
$$

by Lemma 1.6.6. Since $I-(P-\lambda) R(\lambda) \in S^{-n_{0}}$, we use Lemma 1.7.1 to complete the proof.

We define $E(t)=\frac{1}{2 \pi i} \int_{\gamma} e^{-t \lambda} R(\lambda) d \lambda$. We will show shortly that this has a smooth kernel $K^{\prime}(t, x, y)$. Let $K(t, x, y)$ be the smooth kernel of $e^{-t P}$. We
will use Lemma 1.2.9 to estimate the difference between these two kernels. We compute:

$$
E(t)-e^{-t P}=\frac{1}{2 \pi i} \int_{\gamma} e^{-t \lambda}\left(R(\lambda)-(P-\lambda)^{-1}\right) d \lambda
$$

We assume $0<t<1$ and make a change of variables to replace $t \lambda$ by $\lambda$. We use Cauchy's theorem to shift the resulting path $t \gamma$ inside $\mathcal{R}$ back to the original path $\gamma$ where we have uniform estimates. This expresses:

$$
E(t)-e^{-t P}=\frac{1}{2 \pi i} \int_{\gamma} e^{-\lambda}\left(R\left(t^{-1} \lambda\right)-\left(P-t^{-1} \lambda\right)\right) t^{-1} d \lambda
$$

We estimate therefore:

$$
\begin{aligned}
\left|E(t)-e^{-t P}\right|_{-k, k} & \leq C_{k} \int_{\gamma}\left|e^{-\lambda}\left(1+t^{-1}|\lambda|\right)^{-k} d \lambda\right| \\
& \leq C_{k} t^{k}
\end{aligned}
$$

provided $n_{0}$ is large enough. Lemma 1.2.9 implies
Lemma 1.7.3. Let $k$ be given. If $n_{0}$ is large enough, we can estimate:

$$
\left|K(t, x, y)-K^{\prime}(t, x, y)\right|_{\infty, k} \leq C_{k} t^{k} \quad \text { for } 0<t<1
$$

This implies that $K^{\prime}$ approximates $K$ to arbitrarily high jets as $t \rightarrow 0$.
We now study the operator $E(t)$. We define

$$
e_{n}(t, x, \xi)=\frac{1}{2 \pi i} \int_{\gamma} e^{-t \lambda} r_{n}(x, \xi, \lambda) d \lambda
$$

then $E(t)$ is a $\Psi D O$ with symbol $e_{0}+\cdots+e_{n_{0}}$. If we integrate by parts in $\lambda$ we see that we can express

$$
e_{n}(t, x, \xi)=\frac{1}{2 \pi i} \frac{1}{t^{k}} \int_{\gamma} e^{-t \lambda} \frac{d^{k}}{d \lambda^{k}} r_{n}(x, \xi, \lambda) d \lambda
$$

Since $\frac{d^{k}}{d \lambda^{k}} r_{n}$ is homogeneous of degree $-d-n-k d$ in $(\xi, \lambda)$, we see that $e_{n}(t, x, \xi) \in S^{-\infty}$ for any $t>0$. Therefore we can apply Lemma 1.2 .5 to conclude $E_{n}(t)$ which has symbol $e_{n}(t)$ is represented by a kernel function defined by:

$$
K_{n}(t, x, y)=\int e^{i(x-y) \cdot \xi} e_{n}(t, x, \xi) d \xi
$$

We compute:

$$
K_{n}(t, x, x)=\frac{1}{2 \pi i} \iint_{\gamma} e^{-t \lambda} r_{n}(x, \xi, \lambda) d \lambda d \xi
$$

We make a change of variables to replace $\lambda$ by $t^{-1} \lambda$ and $\xi$ by $t^{-1 / d} \xi$ to compute:

$$
\begin{aligned}
K_{n}(t, x, x) & =t^{-\frac{m}{d}-1} \frac{1}{2 \pi i} \iint_{\gamma} e^{-\lambda} r_{n}\left(x, t^{\frac{-1}{d}} \xi, t^{-1} \lambda\right) d \lambda d \xi \\
& =t^{-\frac{m}{d}-1+\frac{n+d}{d}} \frac{1}{2 \pi i} \iint_{\gamma} e^{-\lambda} r_{n}(x, \xi, \lambda) d \lambda d \xi \\
& =t^{\frac{n-m}{d}} e_{n}(x)
\end{aligned}
$$

where we let this integral define $e_{n}(x)$.
Since $p_{d}(x, \xi)$ is assumed to be positive definite, $d$ must be even since $p_{d}$ is a polynomial in $\xi$. Inductively we express $r_{n}$ as a sum of terms of the form:

$$
r_{0}^{j_{1}} q_{1} r_{0}^{j_{2}} q_{2} \ldots r_{0}^{j_{k}}
$$

where the $q_{k}$ are polynomials in $(x, \xi)$. The sum of the degrees of the $q_{k}$ is odd if $n$ is odd and therefore $r_{n}(x,-\xi, \lambda)=-r_{n}(x, \xi, \lambda)$. If we replace $\xi$ by $-\xi$ in the integral defining $e_{n}(x)$, we conclude $e_{n}(x)=0$ if $n$ is odd.

Lemma 1.7.4. Let $P$ be a self-adjoint elliptic partial differential operator of order $d>0$ such that the leading symbol of $P$ is positive definite for $\xi \neq 0$. Then:
(a) If we choose a coordinate system for $M$ near a point $x \in M$ and choose a local frame for $V$, we can define $e_{n}(x)$ using the complicated combinatorial recipe given above. $e_{n}(x)$ depends functorially on a finite number of jets of the symbol $p(x, \xi)$.
(b) If $K(t, x, y)$ is the kernel of $e^{-t P}$ then

$$
K(t, x, x) \sim \sum_{n=0}^{\infty} t^{\frac{n-m}{d}} e_{n}(x) \quad \text { as } t \rightarrow 0^{+}
$$

i.e., given any integer $k$ there exists $n(k)$ such that:

$$
\left|K(t, x, x)-\sum_{n \leq n(k)} t^{\frac{n-m}{d}} e_{n}(x)\right|_{\infty, k}<C_{k} t^{k} \quad \text { for } 0<t<1
$$

(c) $e_{n}(x) \in \operatorname{END}(V, V)$ is invariantly defined independent of the coordinate system and local frame for $V$.
(d) $e_{n}(x)=0$ if $n$ is odd.

Proof: (a) is immediate. We computed that

$$
K^{\prime}(t, x, x)=\sum_{n=0}^{n_{0}} t^{\frac{n-m}{d}} e_{n}(x)
$$

so (b) follows from Lemma 1.7.3. Since $K(t, x, x)$ does not depend on the coordinate system chosen, (c) follows from (b). We computed (d) earlier to complete the proof.

We remark that this asymptotic representation of $K(t, x, x)$ exists for a much wider class of operators $P$. We refer to the literature for further details. We shall give explicit formulas for $e_{0}, e_{2}$, and $e_{4}$ in section 4.8 for certain examples arising in geometry.

The invariants $e_{n}(x)=e_{n}(x, P)$ are sections to the bundle of endomorphisms, $\operatorname{END}(V)$. They have a number of functorial properties. We sumarize some of these below.

Lemma 1.7.5.
(a) Let $P_{i}: C^{\infty}\left(V_{i}\right) \rightarrow C^{\infty}\left(V_{i}\right)$ be elliptic self-adjoint partial differential operators of order $d>0$ with positive definite leading symbol. We form $P=P_{1} \oplus P_{2}: C^{\infty}\left(V_{1} \oplus V_{2}\right) \rightarrow C^{\infty}\left(V_{1} \oplus V_{2}\right)$. Then $P$ is an elliptic selfadjoint partial differential operator of order $d>0$ with positive definite leading symbol and $e_{n}\left(x, P_{1} \oplus P_{2}\right)=e_{n}\left(x, P_{1}\right) \oplus e_{n}\left(x, P_{2}\right)$.
(b) Let $P_{i}: C^{\infty}\left(V_{i}\right) \rightarrow C^{\infty}\left(V_{i}\right)$ be elliptic self-adjoint partial differential operators of order $d>0$ with positive definite leading symbol defined over different manifolds $M_{i}$. We let

$$
P=P_{1} \otimes 1+1 \otimes P_{2}: C^{\infty}\left(V_{1} \otimes V_{2}\right) \rightarrow C^{\infty}\left(V_{1} \otimes V_{2}\right)
$$

over $M=M_{1} \times M_{2}$. Then $P$ is an elliptic self-adjoint partial differential operator of order $d>0$ with positive definite leading symbol over $M$ and

$$
e_{n}(x, P)=\sum_{p+q=n} e_{p}\left(x_{1}, P_{1}\right) \otimes e_{q}\left(x_{2}, P_{2}\right) .
$$

(c) Let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be an elliptic self-adjoint partial differential operator of order $d>0$ with positive definite leading symbol. We decompose the total symbol of $P$ in the form:

$$
p(x, \xi)=\sum_{|\alpha| \leq d} p_{\alpha}(x) \xi^{\alpha}
$$

Fix a local frame for $V$ and a system of local coordinates on $M$. We let indices $a, b$ index the local frame and let $p_{\alpha}=p_{a b \alpha}$ give the components
of the matrix $p_{\alpha}$. We introduce formal variables $p_{a b \alpha / \beta}=D_{x}^{\beta} p_{a b \alpha}$ for the jets of the symbol of $P$. Define $\operatorname{ord}\left(p_{a b \alpha / \beta}\right)=|\beta|+d-|\alpha|$. Then $e_{n}(x, P)$ can be expressed as a sum of monomials in the $\left\{p_{a b \alpha / \beta}\right\}$ variables which are homogeneous of order $n$ in the jets of the symbol as discussed above and with coefficients which depend smoothly on the leading symbol $\left\{p_{a b \alpha}\right\}_{|\alpha|=d}$.
(d) If the leading symbol of $P$ is scalar, then the invariance theory can be simplified. We do not need to introduce the components of the symbol explicitly and can compute $e_{n}(x, P)$ as a non-commutative polynomial which is homogeneous of order $n$ in the $\left\{p_{\alpha / \beta}\right\}$ variables with coefficients which depend smoothly upon the leading symbol.
Remark: The statement of this lemma is somewhat technical. It will be convenient, however, to have these functorial properties precisely stated for later reference. This result suffices to study the index theorem. We shall give one more result which generalizes this one at the end of this section useful in studying the eta invariant. The objects we are studying have a bigrading; one grading comes from counting the number of derivatives in the jets of the symbol, while the other measures the degree of homogeneity in the $(\xi, \lambda)$ variables.
Proof: (a) and (b) follow from the identities:

$$
\begin{aligned}
e^{-t\left(P_{1} \oplus P_{2}\right)} & =e^{-t P_{1}} \oplus e^{-t P_{2}} \\
e^{-t\left(P_{1} \otimes 1+1 \otimes P_{2}\right)} & =e^{-t P_{1}} \otimes e^{-t P_{2}}
\end{aligned}
$$

so the kernels satisfy the identities:

$$
\begin{aligned}
K\left(t, x, x, P_{1} \oplus P_{2}\right) & =K\left(t, x, x, P_{1}\right) \oplus K\left(t, x, x, P_{2}\right) \\
K\left(t, x, x, P_{1} \otimes 1+1 \otimes P_{2}\right) & =K\left(t, x_{1}, x_{1}, P_{1}\right) \otimes K\left(t, x_{2}, x_{2}, P_{2}\right)
\end{aligned}
$$

We equate equal powers of $t$ in the asymptotic series:

$$
\begin{aligned}
\sum t^{\frac{n-m}{d}} e_{n}(x, & \left.P_{1} \oplus P_{2}\right) \\
& \sim \sum t^{\frac{n-m}{d}} e_{n}\left(x, P_{1}\right) \oplus \sum t^{\frac{n-m}{d}} e_{n}\left(x, P_{2}\right) \\
\sum t^{\frac{n-m}{d}} e_{n}(x, & \left.P_{1} \otimes 1+1 \otimes P_{2}\right) \\
& \sim\left\{\sum t^{\frac{p-m_{1}}{d}} e_{p}\left(x_{1}, P_{1}\right)\right\} \otimes\left\{\sum t^{\frac{q-m_{2}}{d}} e_{q}\left(x_{2}, P_{2}\right)\right\}
\end{aligned}
$$

to complete the proof of (a) and (b). We note that the multiplicative property (b) is a direct consequence of the identity $e^{-t(a+b)}=e^{-t a} e^{-t b}$; it was for this reason we worked with the heat equation. Had we worked instead with the zeta function to study $\operatorname{Tr}\left(P^{-s}\right)$ the corresponding multiplicative property would have been much more difficult to derive.

We prove (c) as follows: expand $p(x, \xi)=\sum_{j} p_{j}(x, \xi)$ into homogeneous polynomials where $p_{j}(x, \xi)=\sum_{|\alpha|=j} p_{\alpha}(x) \xi^{\alpha}$. Suppose for the moment that $p_{d}(x, \xi)=p_{d}(x, \xi) I_{V}$ is scalar so it commutes with every matrix. The approximate resolvant is given by:

$$
\begin{aligned}
& r_{0}=\left(p_{d}(x, \xi)-\lambda\right)^{-1} \\
& r_{n}=-r_{0} \sum_{\substack{|\alpha|+j+d-k=n \\
j<n}} d_{\xi}^{\alpha} r_{j} D_{x}^{\alpha} p_{k} / \alpha!
\end{aligned}
$$

Since $p_{d}$ is scalar, $r_{0}$ is scalar so it and all its derivatives commute with any matrix. If we assume inductively $r_{j}$ is of order $j$ in the jets of the symbol, then $d_{\xi}^{\alpha} r_{j} D_{x}^{\alpha} p_{k} / \alpha$ ! is homogeneous of order $j+|\alpha|+d-k=n$ in the jets of the symbol. We observe inductively we can decompose $r_{n}$ in the form:

$$
r_{n}=\sum_{\substack{n=d j+d-|\alpha| \\|\alpha| \leq n}} r_{0}^{j} r_{n, j, \alpha}(x) \xi^{\alpha}
$$

where the $r_{n, j, \alpha}(x)$ are certain non-commutative polynomials in the jets of the total symbol of $P$ which are homogeneous of order $n$ in the sense we have defined.

The next step in the proof of Lemma 1.7.4 was to define:

$$
\begin{aligned}
e_{n}(t, x, \xi) & =\frac{1}{2 \pi i} \int_{\gamma} e^{-t \lambda} r_{n}(x, \xi, \lambda) d \lambda \\
& =\frac{1}{2 \pi i} \sum_{j, \alpha} r_{n, j, \alpha}(x) \xi^{\alpha} \int_{\gamma} e^{-t \lambda} r_{0}^{j} d \lambda \\
e_{n}(x) & =\frac{1}{2 \pi i} \sum_{j, \alpha} r_{n, j, \alpha}(x) \int\left(\int_{\gamma} e^{-\lambda} r_{0}^{j} d \lambda\right) \xi^{\alpha} d \xi .
\end{aligned}
$$

We note again that $e_{n}=0$ if $n$ is odd since the resulting function of $\xi$ would be odd. The remaining coefficients of $r_{n, j, \alpha}(x)$ depend smoothly on the leading symbol $p_{d}$. This completes the proof of the lemma.

If the leading symbol is not scalar, then the situation is more complicated. Choose a local frame to represent the $p_{\alpha}=p_{a b \alpha}$ as matrices. Let $h(x, \xi, \lambda)=\operatorname{det}\left(p_{d}(x, \xi)-\lambda\right)^{-1}$. By Cramer's rule, we can express $r_{0}=r_{0 a b}=\left\{\left(p_{d}-\lambda\right)^{-1}\right\}_{a b}$ as polynomials in the $\left\{h, \lambda, \xi, p_{a b \alpha}(x)\right\}$ variables. We noted previously that $r_{n}$ was a sum of terms of the form $r_{0} q_{0} r_{0} q_{1} \ldots r_{0} q_{k} r_{0}$. The matrix components of such a product can in turn be decomposed as a sum of terms $h^{v} \tilde{q}_{v}$ where $\tilde{q}_{v}$ is a polynomial in the $\left(\lambda, \xi, p_{a b \alpha / \beta}\right)$ variables. The same induction argument which was used in the scalar case shows the $\tilde{q}_{v}$ will be homogeneous of order $n$ in the jets of
the symbol. The remainder of the argument is the same; performing the $d \lambda d \xi$ integral yields a smooth function of the leading symbol as a coefficient of such a term, but this function is not in general rational of course. The scalar case is much simpler as we don't need to introduce the components of the matrices explicitly (although the frame dependence is still there of course) since $r_{0}$ can be commuted. This is a technical point, but one often useful in making specific calculations.

We define the scalar invariants

$$
a_{n}(x, P)=\operatorname{Tr} e_{n}(x, P)
$$

where the trace is the fiber trace in $V$ over the point $x$. These scalar invariants $a_{n}(x, P)$ inherit suitable functorial properties from the functorial properties of the invariants $e_{n}(x, P)$. It is immediate that:

$$
\begin{aligned}
\operatorname{Tr}_{L^{2}} e^{-t P} & =\int_{M} \operatorname{Tr}_{V_{x}} K(t, x, x) \mathrm{dvol}(x) \\
& \sim \sum_{n=0}^{\infty} t^{\frac{m-n}{d}} \int_{M} a_{n}(x, P) \operatorname{dvol}(x) \\
& \sim \sum_{n=0}^{\infty} t^{\frac{m-n}{d}} a_{n}(P)
\end{aligned}
$$

where $a_{n}(P)=\int_{M} a_{n}(x, P) \operatorname{dvol}(x)$ is the integrated invariant. This is a spectral invariant of $P$ which can be computed from local information about the symbol of $P$.

Let $(P, V)$ be an elliptic complex of differential operators and let $\Delta_{i}$ be the associated Laplacians. We define:

$$
a_{n}(x, P)=\sum_{i}(-1)^{i} \operatorname{Tr} e_{n}\left(x, \Delta_{i}\right)
$$

then Lemma 1.6.5 and the remark which follows this lemma imply

$$
\operatorname{index}(P)=\sum_{i}(-1)^{i} \operatorname{Tr} e^{-t \Delta_{i}} \sim \sum_{n=0}^{\infty} t^{\frac{n-m}{d}} \int_{M} a_{n}(x, P) \operatorname{dvol}(x)
$$

Since the left hand side does not depend on the parameter $t$, we conclude:
Theorem 1.7.6. Let $(P, V)$ be a elliptic complex of differential operators.
(a) $a_{n}(x, P)$ can be computed in any coordinate system and relative to any local frames as a complicated combinatorial expression in the jets of $P$ and of $P^{*}$ up to some finite order. $a_{n}=0$ if $n$ is odd.
(b)

$$
\int_{M} a_{n}(x, P) \operatorname{dvol}(x)= \begin{cases}\operatorname{index}(P) & \text { if } n=m \\ 0 & \text { if } n \neq m\end{cases}
$$

We note as an immediate consequence that $\operatorname{index}(P)=0$ if $m$ is odd. The local nature of the invariants $a_{n}$ will play a very important role in our discussion of the index theorem. We will develop some of their functorial properties at that point. We give one simple example to illustrate this fact.

Let $\pi: M_{1} \rightarrow M_{2}$ be a finite covering projection with fiber $F$. We can choose the metric on $M_{1}$ to be the pull-back of a metric on $M_{2}$ so that $\pi$ is an isometry. Let $d$ be the operator of the de Rham complex. $a_{n}(x, d)=$ $a_{n}(\pi x, d)$ since $a_{n}$ is locally defined. This implies:

$$
\chi\left(M_{1}\right)=\int_{M_{1}} a_{m}(x, d) \operatorname{dvol}(x)=|F| \int_{M_{2}} a_{m}(x, d) \operatorname{dvol}(x)=|F| \chi\left(M_{2}\right)
$$

so the Euler-Poincare characteristic is multiplicative under finite coverings. A similar argument shows the signature is multiplicative under orientation preserving coverings and that the arithmetic genus is multiplicative under holomorphic coverings. (We will discuss the signature and arithmetic genus in more detail in Chapter 3).

We conclude this section with a minor generalization of Lemma 1.7.5 which will be useful in discussing the eta invariant in section 1.10:

Lemma 1.7.7. Let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be an elliptic self-adjoint partial differential operator of order $d>0$ with positive definite leading symbol and let $Q$ be an auxilary partial differential operator on $C^{\infty}(V)$ of order $a \geq 0 . Q e^{-t P}$ is an infinitely smoothing operator with kernel $Q K(t, x, y)$. There is an asymptotic expansion on the diagonal:

$$
\left.\{Q K(t, x, y)\}\right|_{x=y} \sim \sum_{n=0}^{\infty} t^{(n-m-a) / d} e_{n}(x, Q, P)
$$

The $e_{n}$ are smooth local invariants of the jets of the symbols of $P$ and $Q$ and $e_{n}=0$ if $n+a$ is odd. If we let $Q=\sum q_{\alpha} D_{x}^{\alpha}$, we define $\operatorname{ord}\left(q_{\alpha}\right)=a-|\alpha|$. If the leading symbol of $P$ is scalar, we can compute $e_{n}(x, Q, P)$ as a non-commutative polynomial in the variables $\left\{q_{\alpha}, p_{\alpha / \beta}\right\}$ which is homogeneous of order $n$ in the jets with coefficients which depend smoothly on the $\left\{p_{\alpha}\right\}_{|\alpha|=d}$ variables. This expression is linear in the $\left\{q_{\alpha}\right\}$ variables and does not involve the higher jets of these variables. If the leading symbol of $P$ is not scalar, there is a similar expression for the matrix components of $e_{n}$ in the matrix components of these variables. $e_{n}(x, Q, P)$ is additive and multiplicative in the sense of Lemma 1.7.5(a) and (b) with respect to direct sums of operators and tensor products over product manifolds. Remark: In fact it is not necessary to assume $P$ is self-adjoint to define $e^{-t P}$ and to define the asymptotic series. It is easy to generalize the techniques
we have developed to prove this lemma continues to hold true if we only asume that $\operatorname{det}\left(p_{d}(x, \xi)-\lambda\right) \neq 0$ for $\xi \neq 0$ and $\operatorname{Im}(\lambda) \leq 0$. This implies that the spectrum of $P$ is pure point and contained in a cone about the positive real axis. We omit the details.
Proof: $Q e^{-t P}$ has smooth kernel $Q K(t, x, y)$ where $Q$ acts as a differential operator on the $x$ variables. One representative piece is:

$$
\begin{aligned}
q_{\alpha}(x) D_{x}^{\alpha} K_{n}(t, x, y) & =q_{\alpha}(x) D_{x}^{\alpha} \int e^{i(x-y) \cdot \xi} e_{n}(t, x, \xi) d \xi \\
& =q_{\alpha}(x) \sum_{\beta+\gamma=\alpha} \frac{\alpha!}{\beta!\gamma!} \int e^{i(x-y) \cdot \xi} \xi^{\beta} D_{x}^{\gamma} e_{n}(t, x, \xi) d \xi
\end{aligned}
$$

We define $q_{\beta, \gamma}(x)=q_{\alpha}(x) \frac{\alpha!}{\beta!\gamma!}$ where $\beta+\gamma=\alpha$, then a representative term of this kernel has the form:

$$
q_{\beta, \gamma}(x) \int e^{i(x-y) \cdot \xi} \xi^{\beta} D_{x}^{\gamma} e_{n}(t, x, \xi) d \xi
$$

where $q_{\beta, \gamma}(x)$ is homogeneous of order $a-|\beta|-|\gamma|$ in the jets of the symbol of $Q$.

We suppose the leading symbol is scalar to simplify the computations; the general case is handled similarly using Cramer's rule. We express

$$
e_{n}=\frac{1}{2 \pi i} \int e^{-t \lambda} r_{n}(x, \xi, \lambda) d \lambda
$$

where $r_{n}$ is a sum of terms of the form $r_{n, j, \alpha}(x) \xi^{\alpha} r_{0}^{j}$. When we differentiate such a term with respect to the $x$ variables, we change the order in the jets of the symbol, but do not change the $(\xi, \lambda)$ degree of homogeneity. Therefore $Q K(t, x, y)$ is a sum of terms of the form:

$$
q_{\beta, \gamma}(x) \int\left\{\int e^{-t \lambda} r_{0}^{j}(x, \xi, \lambda) d \lambda\right\} e^{i(x-y) \cdot \xi} \xi^{\beta} \xi^{\delta} \tilde{r}_{n, j, \delta}(x) d \xi
$$

where:

$$
\begin{aligned}
& r_{n, j, \delta} \text { is of order } n+|\gamma| \text { in the jets of the symbol, } \\
& -d-n=-j d+|\delta|
\end{aligned}
$$

We evaluate on the diagonal to set $x=y$. This expression is homogeneous of order $-d-n+|\beta|$ in $(\xi, \lambda)$, so when we make the appropriate change of variables we calculate this term becomes:

$$
t^{(n-m-|\beta|) / d} q_{\beta, \gamma}(x) \tilde{r}_{n, j, \delta}(x) \cdot\left\{\int\left\{e^{-\lambda} r_{0}^{j}(x, \xi, \lambda) d \lambda\right\} \xi^{\beta} \xi^{\delta} d \xi\right\}
$$

This is of order $n+|\gamma|+a-|\beta|-|\gamma|=n+a-|\beta|=\nu$ in the jets of the symbol. The exponent of $t$ is $(n-m-|\beta|) / d=(\nu-a+|\beta|-m-|\beta|) / d=(\nu-a-m) / d$. The asymptotic formula has the proper form if we index by the order in the jets of the symbol $\nu$. It is linear in $q$ and vanishes if $|\beta|+|\delta|$ is odd. As $d$ is even, we compute:

$$
\nu+a=n+a-|\beta|+a=j d-d-|\delta|+2 a-|\beta| \equiv|\delta|+|\beta| \quad \bmod 2
$$

so that this term vanishes if $\nu+a$ is odd. This completes the proof.

### 1.8. Lefschetz Fixed Point Theorems.

Let $T: M \rightarrow M$ be a continuous map. $T^{*}$ defines a map on the cohomology of $M$ and the Lefschetz number of $T$ is defined by:

$$
L(T)=\sum(-1)^{p} \operatorname{Tr}\left(T^{*} \text { on } H^{p}(M ; \mathbf{C})\right)
$$

This is always an integer. We present the following example to illustrate the concepts involved. Let $M=S^{1} \times S^{1}$ be the two dimensional torus which we realize as $\mathbf{R}^{2}$ modulo the integer lattice $\mathbf{Z}^{2}$. We define $T(x, y)=$ $\left(n_{1} x+n_{2} y, n_{3} x+n_{4} y\right)$ where the $n_{i}$ are integers. Since this preserves the integer lattice, this defines a map from $M$ to itself. We compute:

$$
\begin{aligned}
& T^{*}(1)=1, \quad T^{*}(d x \wedge d y)=\left(n_{1} n_{4}-n_{2} n_{3}\right) d x \wedge d y \\
& T^{*}(d x)=n_{1} d x+n_{2} d y, \quad T^{*}(d y)=n_{3} d x+n_{4} d y \\
& L(T)=1+\left(n_{1} n_{4}-n_{2} n_{3}\right)-\left(n_{1}+n_{4}\right)
\end{aligned}
$$

Of course, there are many other interesting examples.
We computed $L(T)$ in the above example using the de Rham isomorphism. We let $T^{*}=\Lambda^{p}(d T): \Lambda^{p}\left(T^{*} M\right) \rightarrow \Lambda^{p}\left(T^{*} M\right)$ to be the pull-back operation. It is a map from the fiber over $T(x)$ to the fiber over $x$. Since $d T^{*}=T^{*} d, T^{*}$ induces a map on $H^{p}(M, \mathbf{C})=\operatorname{ker}\left(d_{p}\right) / \operatorname{image}\left(d_{p-1}\right)$. If $T$ is the identity map, then $L(T)=\chi(M)$. The perhaps somewhat surprising fact is that Lemma 1.6 .5 can be generalized to compute $L(T)$ in terms of the heat equation.

Let $\Delta_{p}=(d \delta+\delta d)_{p}: C^{\infty}\left(\Lambda^{p} T^{*} M\right) \rightarrow C^{\infty}\left(\Lambda^{p} T^{*} M\right)$ be the associated Laplacian. We decompose $L^{2}\left(\Lambda^{p} T^{*} M\right)=\bigoplus_{\lambda} E_{p}(\lambda)$ into the eigenspaces of $\Delta_{p}$. We let $\pi(p, \lambda)$ denote orthogonal projection on these subspaces, and we define $T^{*}(p, \lambda)=\pi(p, \lambda) T^{*}: E_{p}(\lambda) \rightarrow E_{p}(\lambda)$. It is immediate that:

$$
\operatorname{Tr}\left(T^{*} e^{-t \Delta_{p}}\right)=\sum_{\lambda} e^{-t \lambda} \operatorname{Tr}\left(T^{*}(p, \lambda)\right)
$$

Since $d T^{*}=T^{*} d$ and $d \pi=\pi d$, for $\lambda \neq 0$ we get a chain map between the exact sequences:


Since this diagram commutes and since the two rows are long exact sequences of finite dimensional vector spaces, a standard result in homological algebra implies:

$$
\sum_{p}(-1)^{p} \operatorname{Tr}\left(T^{*}(p, \lambda)\right)=0 \quad \text { for } \lambda \neq 0
$$

It is easy to see the corresponding sum for $\lambda=0$ yields $L(T)$, so Lemma 1.6.5 generalizes to give a heat equation formula for $L(T)$.

This computation was purely formal and did not depend on the fact that we were dealing with the de Rham complex.

Lemma 1.8.1. Let $(P, V)$ be an elliptic complex over $M$ and let $T: M \rightarrow$ $M$ be smooth. We assume given linear maps $V_{i}(T): V_{i}(T x) \rightarrow V_{i}(x)$ so that $P_{i}\left(V_{i}(T)\right)=V_{i}(T) P_{i}$. Then $T$ induces a map on $H^{p}(P, V)$. We define

$$
L(T)_{P}=\sum_{p}(-1)^{p} \operatorname{Tr}\left(T \text { on } H^{p}(P, V)\right) .
$$

Then we can compute $L(T)_{P}=\sum_{p} \operatorname{Tr}\left(V_{i}(T) e^{-t \Delta_{p}}\right)$ where $\Delta_{p}=\left(P^{*} P+\right.$ $\left.P P^{*}\right)_{p}$ is the associated Laplacian.

The $V_{i}(T)$ are a smooth linear action of $T$ on the bundles $V_{i}$. They will be given by the representations involved for the de Rham, signature, spin, and Dolbeault complexes as we shall discuss in Chapter 4. We usually denote $V_{i}(T)$ by $T^{*}$ unless it is necessary to specify the action involved.

This Lemma implies Lemma 1.6.5 if we take $T=I$ and $V_{i}(T)=I$. To generalize Lemma 1.7.4 and thereby get a local formula for $L(T)_{P}$, we must place some restrictions on the map $T$. We assume the fixed point set of $T$ consists of the finite disjoint union of smooth submanifolds $N_{i}$. Let $\nu\left(N_{i}\right)=T(M) / T\left(N_{i}\right)$ be the normal bundle over the submanifold $N_{i}$. Since $d T$ preserves $T\left(N_{i}\right)$, it induces a map $d T_{\nu}$ on the bundle $\nu\left(N_{i}\right)$. We suppose $\operatorname{det}\left(I-d T_{\nu}\right) \neq 0$ as a non-degeneracy condition; there are no normal directions left fixed infinitesimally by $T$.

If $T$ is an isometry, this condition is automatic. We can construct a non-example by defining $T(z)=z /(z+1): S^{2} \rightarrow S^{2}$. The only fixed point is at $z=0$ and $d T(0)=I$, so this fixed point is degenerate.

If $K$ is the kernel of $e^{-t P}$, we pull back the kernel to define $T^{*}(K)(t, x, y)$ $=T^{*}(x) K(t, T x, y)$. It is immediate $T^{*} K$ is the kernel of $T^{*} e^{-t P}$.

Lemma 1.8.2. Let $P$ be an elliptic partial differential operator of order $d>0$ which is self-adjoint and which has a positive definite leading symbol for $\xi \neq 0$. Let $T: M \rightarrow M$ be a smooth non-degenerate map and let $T^{*}: V_{T(x)} \rightarrow V_{x}$ be a smooth linear action. If $K$ is the kernel of $e^{-t P}$ then $\operatorname{Tr}\left(T^{*} e^{-t P}\right)=\int_{M} \operatorname{Tr}\left(T^{*} K\right)(t, x, x) \mathrm{dvol}(x)$. Furthermore:
(a) If $T$ has no fixed points, $\left|\operatorname{Tr}\left(T^{*} e^{-t P}\right)\right| \leq C_{n} t^{-n}$ as $t \rightarrow 0^{+}$for any $n$.
(b) If the fixed point set of $T$ consists of submanifolds $N_{i}$ of dimension $m_{i}$, we will construct scalar invariants $a_{n}(x)$ which depend functorially upon a finite number of jets of the symbol and of $T$. The $a_{n}(x)$ are defined over $N_{i}$ and

$$
\operatorname{Tr}\left(T^{*} e^{-t P}\right) \sim \sum_{i} \sum_{n=0}^{\infty} t^{\frac{n-m_{i}}{d}} \int_{N_{i}} a_{n}(x) \operatorname{dvol}_{i}(x)
$$

$\operatorname{dvol}_{i}(x)$ denotes the Riemannian measure on the submanifold.
It follows that if $T$ has no fixed points, then $L(T)_{P}=0$.
Proof: Let $\left\{e_{n}, r_{n}, K_{n}\right\}$ be as defined in section 1.7. The estimates of Lemma 1.7 .3 show that $\left|T^{*} K-\sum_{n \leq n_{0}} T^{*} K_{n}\right|_{\infty, k} \leq C(k) t^{-k}$ as $t \rightarrow 0^{+}$ for any $k$ if $n_{0}=n_{0}(k)$. We may therefore replace $K$ by $K_{n}$ in proving (a) and (b). We recall that:

$$
e_{n}(t, x, \xi)=\frac{1}{2 \pi i} \int_{\gamma} e^{-t \lambda} r_{n}(x, \xi, \lambda) d \lambda
$$

We use the homogeneity of $r_{n}$ to express:

$$
\begin{aligned}
e_{n}(t, x, \xi) & =\frac{1}{2 \pi i} \int_{\gamma} e^{-\lambda} r_{n}\left(x, \xi, t^{-1} \lambda\right) t^{-1} d \lambda \\
& =t^{\frac{n}{d}} \frac{1}{2 \pi i} \int_{\gamma} e^{-\lambda} r_{n}\left(x, t^{\frac{1}{d}} \xi, \lambda\right) d \lambda \\
& =t^{\frac{n}{d}} e_{n}\left(x, t^{\frac{1}{d}} \xi\right)
\end{aligned}
$$

where we define $e_{n}(x, \xi)=e_{n}(1, x, \xi) \in S^{-\infty}$. Then:

$$
\begin{aligned}
K_{n}(t, x, y) & =\int e^{i(x-y) \cdot \xi} e_{n}(t, x, \xi) d \xi \\
& =t^{\frac{n-m}{d}} \int e^{i(x-y) \cdot t^{-\frac{1}{d} \xi}} e_{n}(x, \xi) d \xi
\end{aligned}
$$

This shows that:

$$
T^{*} K_{n}(t, x, x)=t^{\frac{n-m}{d}} \int T^{*}(x) e^{i(T x-x) \cdot t^{-1 / d} \xi} e_{n}(T x, \xi) d \xi
$$

We must study terms which have the form:

$$
\int e^{i(T x-x) \cdot t^{-1 / d} \xi} \operatorname{Tr}\left(T^{*}(x) e_{n}(T x, \xi)\right) d \xi d x
$$

where the integral is over the cotangent space $T^{*}(M)$. We use the method of stationary phase on this highly oscillatory integral. We first bound
$|T x-x| \geq \varepsilon>0$. Using the argument developed in Lemma 1.2.6, we integrate by parts to bound this integral by $C(n, k, \varepsilon) t^{k}$ as $t \rightarrow 0$ for any $k$. If $T$ has no fixed points, this proves (a). There is a slight amount of notational sloppiness here since we really should introduce partitions of unity and coordinate charts to define $T x-x$, but we supress these details in the interests of clarity.

We can localize the integral to an arbitrarily small neighborhood of the fixed point set in proving (b). We shall assume for notational simplicity that the fixed point set of $T$ consists of a single submanifold $N$ of dimension $m_{1}$. The map $d T_{\nu}$ on the normal bundle has no eigenvalue 1 . We identify $\nu$ with the span of the generalized eigenvectors of $d T$ on $\left.T(M)\right|_{N}$ which correspond to eigenvalues not equal to 1 . This gives a direct sum decomposition over $N$ :

$$
T(M)=T(N) \oplus \nu \quad \text { and } \quad d T=I \oplus d T_{\nu}
$$

We choose a Riemannian metric for $M$ so this splitting is orthogonal. We emphasize that these are bundles over $N$ and not over the whole manifold $M$.

We describe the geometry near the fixed manifold $N$ using the normal bundle $\nu$. Let $y=\left(y_{1}, \ldots, y_{m_{1}}\right)$ be local coordinates on $N$ and let $\left\{\vec{s}_{1}, \ldots, \vec{s}_{m-m_{1}}\right\}$ be a local orthonormal frame for $\nu$. We use this local orthonormal frame to introduce fiber coordinates $z=\left(z_{1}, \ldots, z_{m-m_{1}}\right)$ for $\nu$ by decomposing any $\vec{s} \in \nu$ in the form:

$$
\vec{s}=\sum_{j} z_{j} \vec{s}_{j}(y)
$$

We let $x=(y, z)$ be local coordinates for $\nu$. The geodesic flow identifies a neighborhood of the zero section of the bundle $\nu$ with a neighborhood of $N$ in $M$ so we can also regard $x=(y, z)$ as local coordinates on $M$.

We decompose

$$
T(x)=\left(T_{1}(x), T_{2}(x)\right)
$$

into tangential and fiber coordinates. Because the Jacobian matrix has the form:

$$
d T(y, 0)=\left(\begin{array}{cc}
I & 0 \\
0 & d T_{\nu}
\end{array}\right)
$$

we conclude that $T_{1}(x)-y$ must vanish to second order in $z$ along $N$.
We integrate $\operatorname{Tr}\left(T^{*} K_{n}\right)(t, x, x)$ along a small neighborhood of the zero section of $\nu$. We shall integrate along the fibers first to reduce this to an integral along $N$. We decompose $\xi=\left(\xi_{1}, \xi_{2}\right)$ corresponding to the decomposition of $x=(y, z)$. Let

$$
D(\nu)=\{(y, z):|z| \leq 1\}
$$

be the unit disk bundle of the normal bundle. Let $U=T^{*}(D(\nu))$ be the cotangent bundle of the unit disk bundle of the normal bundle. We assume the metric chosen so the geodesic flow embeds $D(\nu)$ in $M$. We parametrize $U$ by $\left\{\left(y, z, \xi_{1}, \xi_{2}\right):|z| \leq 1\right\}$. Let $s=t^{\frac{1}{d}}$. Modulo terms which vanish to infinite order in $s$, we compute:

$$
\begin{aligned}
& \mathcal{I} \stackrel{\text { def }}{=} \int_{M} \operatorname{Tr}\left(\left(T^{*} K_{n}\right)(t, x, x)\right) d x \\
&=s^{n-m} \int_{U} e^{i\left(T_{1}(y, z)-y\right) \cdot \xi_{1} s^{-1}} e^{i\left(T_{2}(y, z)-z\right) \cdot \xi_{2} s^{-1}} \\
& \times \operatorname{Tr}\left(T^{*}(x) e_{n}\left(T x, \xi_{1}, \xi_{2}\right)\right) d \xi_{1} d \xi_{2} d z d y .
\end{aligned}
$$

The non-degeneracy assumption on $T$ means the phase function $\bar{w}=$ $T_{2}(y, z)-z$ defines a non-degenerate change of variables if we replace $(y, z)$ by $(y, \bar{w})$. This transforms the integral into the form:

$$
\begin{aligned}
\mathcal{I}= & s^{n-m} \iint_{\bar{U}} e^{i\left(T_{1}(y, \bar{w})-y\right) \cdot \xi_{1} s^{-1}} e^{i \bar{w} \cdot \xi_{2} s^{-1}} \\
& \times\left|\operatorname{det}\left(I-d T_{2}\right)\right|^{-1} \operatorname{Tr}\left(T^{*}(y, \bar{w}) e_{n}\left(T(y, \bar{w}), \xi_{1}, \xi_{2}\right) d \xi_{1} d \xi_{2} d \bar{w} d y\right.
\end{aligned}
$$

where $\bar{U}$ is the image of $U$ under this change. We now make another change of coordinates to let $w=s^{-1} \bar{w}$. As $s \rightarrow 0, s^{-1} \bar{U}$ will converge to $T^{*}(v)$. Since $d \bar{w}=s^{m-m_{1}} d w$ this transforms the integral $I$ to the form:

$$
\begin{aligned}
& \mathcal{I}=s^{n-m_{1}} \int_{s^{-1} \bar{U}} e^{i\left(T_{1}(y, s w)-y\right) \cdot \xi_{1} s^{-1}} e^{i w \cdot \xi_{2}}\left|\operatorname{det}\left(I-d T_{2}\right)^{-1}\right|(y, s w) \\
& \times \operatorname{Tr}\left(T^{*}(y, s w) e_{n}\left(T(y, s w), \xi_{1}, \xi_{2}\right)\right) d \xi_{1} d \xi_{2} d w d y
\end{aligned}
$$

Since the phase function $T_{1}(y, \bar{w})-y$ vanishes to second order at $\bar{w}=0$, the function $\left(T_{1}(y, s w)-y\right) s^{-1}$ is regular at $s=0$. Define:

$$
\begin{aligned}
e_{n}^{\prime}\left(y, w, \xi_{1}, \xi_{2}, s\right)=e^{i\left(T_{1}(y, s w)-y\right) \cdot \xi_{1} s^{-1}} \mid & \operatorname{det}\left(I-d T_{2}\right)(y, s w) \mid \\
\times & \operatorname{Tr}\left(T^{*}(y, s w) e_{n}\left(T(y, s w), \xi_{1}, \xi_{2}\right)\right) .
\end{aligned}
$$

This vanishes to infinite order in $\xi_{1}, \xi_{2}$ at $\infty$ and is regular at $s=0$. To complete the proof of Lemma 1.8.2, we must evaluate:

$$
s^{\left(n-m_{1}\right)} \int_{s^{-1} \bar{U}} e_{n}^{\prime}\left(y, w, \xi_{1}, \xi_{2}, s\right) e^{i w \cdot \xi_{2}} d \xi_{1} d \xi_{2} d w d y
$$

We expand $e_{n}^{\prime}$ in a Taylor series in $s$ centered at $s=0$. If we differentiate $e_{n}^{\prime}$ with respect to $s$ a total of $k$ times and evaluate at $s=0$, then the exponential term disappears and we are left with an expression which is
polynomial in $w$ and of degree at worst $2 k$. It still vanishes to infinite order in $\left(\xi_{1}, \xi_{2}\right)$. We decompose:

$$
e_{n}^{\prime}=\sum_{j \leq j_{0}} s^{j} \sum_{|\alpha| \leq 2 j} c_{j, \alpha}\left(y, \xi_{1}, \xi_{2}\right) w^{\alpha}+s^{j_{0}} \varepsilon\left(y, w, \xi_{1}, \xi_{2}, s\right)
$$

where $\varepsilon$ is the remainder term.
We first study a term of the form:

$$
s^{n+j-m_{1}} \int_{|w|<s^{-1}} c_{j, \alpha}\left(y, \xi_{1}, \xi_{2}\right) w^{\alpha} e^{i \xi_{2} \cdot w} d \xi_{1} d \xi_{2} d w d y
$$

Since $c_{j, \alpha}$ vanishes to infinite order in $\xi$, the $d \xi$ integral creates a function which vanishes to infinite order in $w$. We can let $s \rightarrow 0$ to replace the domain of integration by the normal fiber. The error vanishes to infinite order in $s$ and gives a smooth function of $y$. The $d w$ integral just yields the inverse Fourier transform with appropriate terms and gives rise to asymptotics of the proper form. The error term in the Taylor series grows at worst polynomially in $w$ and can be bounded similarly. This completes the proof. If there is more than one component of the fixed point set, we sum over the components since each component makes a separate contibution to the asymptotic series.

The case of isolated fixed points is of particular interest. Let $d=2$ and $\sigma_{L}(P)=|\xi|^{2} I_{V}$. We compute the first term at a fixed point:

$$
\begin{aligned}
& r_{0}(x, \xi, \lambda)=\left(|\xi|^{2}-\lambda\right)^{-1} \quad e_{0}(x, \xi, t)=e^{-t|\xi|^{2}} \\
& \int \operatorname{Tr}\left(T^{*} K_{0}\right)(t, x, x) \operatorname{dvol}(x)=\int \operatorname{Tr}\left(T^{*}(x)\right) e^{i(T x-x) \cdot \xi} e^{-t|\xi|^{2}} d \xi d x
\end{aligned}
$$

We assume $T(0)=0$ is an isolated non-degenerate fixed point. We let $y=T x-x$ be a change of variables and compute the first term:

$$
\int \operatorname{Tr}\left(T^{*}(y)\right) e^{i y \cdot \xi}|\operatorname{det}(I-d T(y))|^{-1} e^{-t|\xi|^{2}} d \xi d y
$$

We make a change of variables $\xi \rightarrow \xi t^{-1 / 2}$ and $y \rightarrow y t^{1 / 2}$ to express the first term:

$$
\int \operatorname{Tr}\left(T^{*}\left(y t^{1 / 2}\right)\right) e^{i y \cdot \xi} \mid \operatorname{det}\left(I-d T\left(y t^{1 / 2}\right) \mid e^{-|\xi|^{2}} d \xi d y\right.
$$

The $d \xi$ integral just gives $e^{-|y|^{2}}$ so this becomes

$$
\int \operatorname{Tr}\left(T^{*}\left(y t^{1 / 2}\right)\right)\left|\operatorname{det}\left(I-d T\left(y t^{1 / 2}\right)\right)\right|^{-1} e^{-|y|^{2}} d y
$$

We expand this in a Taylor series at $t=0$ and evaluate to get

$$
\operatorname{Tr} T^{*}(0)|\operatorname{det}(I-d T(O))|^{-1}
$$

This proves:

Lemma 1.8.3. Let $P$ be a second order elliptic partial differential operator with leading symbol $|\xi|^{2} I$. Let $T: M \rightarrow M$ be smooth with nondegenerate isolated fixed points. Then:

$$
\operatorname{Tr} T^{*} e^{-t P}=\sum_{i} \operatorname{Tr}\left(T^{*}\right)|\operatorname{det}(I-d T)|^{-1}\left(x_{i}\right)
$$

summed over the fixed point set.
We combine Lemmas 1.8.1 and 1.8.2 to constuct a local formula for $L(T)_{P}$ to generalize the local formula for index $(P)$ given by Theorem 1.7.6; we will discuss this further in the fourth chapter.

We can use Lemma 1.8 .3 to prove the classical Lefschetz fixed point formula for the de Rham complex. Let $T: V \rightarrow V$ be a linear map, then it is easily computed that:

$$
\sum_{p}(-1)^{p} \operatorname{Tr} \Lambda^{p}(T)=\operatorname{det}(I-T)
$$

We compute:

$$
\begin{aligned}
L(T) & =\sum_{p}(-1)^{p} \operatorname{Tr}\left(T^{*} e^{-t \Delta_{p}}\right) \\
& =\sum_{p, i}(-1)^{p} \operatorname{Tr}\left(\Lambda^{p} d T\right)|\operatorname{det}(I-d T)|^{-1}\left(x_{i}\right) \\
& =\sum_{i} \operatorname{det}(I-d T)|\operatorname{det}(I-d T)|^{-1}\left(x_{i}\right) \\
& =\sum_{i} \operatorname{sign} \operatorname{det}(I-d T)\left(x_{i}\right)
\end{aligned}
$$

summed over the fixed point set of $T$. This proves:
Theorem 1.8.4 (Classical Lefschetz Fixed Point Formula).
Let $T: M \rightarrow M$ be smooth with isolated non-degenerate fixed points. Then:

$$
\begin{aligned}
L(T) & =\sum_{p}(-1)^{p} \operatorname{Tr}\left(T^{*} \text { on } H^{p}(M ; \mathbf{C})\right) \\
& =\sum_{i} \operatorname{sign} \operatorname{det}(I-d T)\left(x_{i}\right)
\end{aligned}
$$

summed over the fixed point set.
Remark: We can generalize Lemma 1.8.2 to study $\operatorname{Tr}\left(T^{*} Q e^{-t P}\right)$ where $Q$ is an auxilary differential operator of order $a$. Just as in lemma 1.7.7 we may obtain an asymptotic series:

$$
\operatorname{Tr}\left(T^{*} Q e^{-t P}\right) \sim \sum_{i} \sum_{n=0}^{\infty} t^{\left(n-m_{i}-a\right) / d} a_{n}(x, Q, P) \operatorname{dvol}_{i}(x)
$$

We shall omit the details as the additional terms created by the operator $Q$ are exactly the same as those given in the proof of Lemma 1.7.7. Each term $a_{n}$ is homogeneous of order $n$ in the jets of the symbols of $(Q, P)$ and of the map $T$ in a suitable sense.

### 1.9. Elliptic Boundary Value Problems.

In section 1.7 we derived a local formula for the index of an elliptic partial-differential complex using heat equation methods. This formula will lead to a heat equation proof of the Atiyah-Singer index theorem which we shall discuss later. Unfortunately, it is not known at present how to give a heat equation proof of the Atiyah-Bott index theorem for manifolds with boundary in full generality. We must adopt a much stronger notion of ellipticity to deal with the analytic problems involved. This will yield a heat equation proof of the Gauss-Bonnet theorem for manifolds with boundary which we shall discuss in the fourth chapter.

Let $M$ be a smooth compact manifold with smooth boundary $d M$ and let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be a partial differential operator of order $d>0$. We let $p=\sigma_{L}(P)$ be the leading symbol of $P$. We assume henceforth that $p$ is self-adjoint and elliptic-i.e., $\operatorname{det} p(x, \xi) \neq 0$ for $\xi \neq 0$. Let $\mathbf{R}_{ \pm}$denote the non-zero positive/negative real numbers. It is immediate that

$$
\operatorname{det}\{p(x, \xi)-\lambda\} \neq 0 \quad \text { for }(\xi, \lambda) \neq(0,0) \in T^{*}(M) \times\left\{\mathbf{C}-\mathbf{R}_{+}-\mathbf{R}_{-}\right\}
$$

since $p$ is self-adjoint and elliptic.
We fix a fiber metric on $V$ and a volume element on $M$ to define the global inner product $(\cdot, \cdot)$ on $L^{2}(V)$. We assume that $P$ is formally selfadjoint:

$$
(P f, g)=(f, P g)
$$

for $f$ and $g$ smooth section with supports disjoint from the boundary $d M$. We must impose boundary conditions to make $P$ self-adjoint. For example, if $P=-\partial^{2} / \partial x^{2}$ on the line segment $[0, A]$, then $P$ is formally self-adjoint and elliptic, but we must impose Neumann or Dirichlet boundary conditions to ensure that $P$ is self-adjoint with discrete spectrum.

Near $d M$ we let $x=(y, r)$ where $y=\left(y_{1}, \ldots, y_{m-1}\right)$ is a system of local coordinates on $d M$ and where $r$ is the normal distance to the boundary. We assume $d M=\{x: r(x)=0\}$ and that $\partial / \partial r$ is the inward unit normal. We further normalize the choice of coordinate by requiring the curves $x(r)=$ $\left(y_{0}, r\right)$ for $r \in[0, \delta)$ are unit speed geodesics for any $y_{0} \in d M$. The inward geodesic flow identifies a neighborhood of $d M$ in $M$ with the collar $d M \times$ $[0, \delta)$ for some $\delta>0$. The collaring gives a splitting of $T(M)=T(d M) \oplus$ $T(\mathbf{R})$ and a dual splitting $T^{*}(M)=T^{*}(d M) \oplus T^{*}(\mathbf{R})$. We let $\xi=(\zeta, z)$ for $\zeta \in T^{*}(d M)$ and $z \in T^{*}(\mathbf{R})$ reflect this splitting.

It is convenient to discuss boundary conditions in the context of graded vector bundles. A graded bundle $U$ over $M$ is a vector bundle $U$ together with a fixed decomposition

$$
U=U_{0} \oplus \cdots \oplus U_{d-1}
$$

into sub-bundles $U_{j}$ where $U_{j}=\{0\}$ is permitted in this decomposition. We let $W=V \otimes 1^{d}=V \oplus \cdots \oplus V$ restricted to $d M$ be the bundle of

Cauchy data. If $W_{j}=\left.V\right|_{d M}$ is the $(j+1)^{\text {st }}$ factor in this decomposition, then this defines a natural grading on $W$; we identify $W_{j}$ with the bundle of $j^{\text {th }}$ normal derivatives. The restriction map:

$$
\underline{\gamma}: C^{\infty}(V) \rightarrow C^{\infty}(W)
$$

defined by:

$$
\underline{\gamma}(f)=\left(f_{0}, \ldots, f_{d-1}\right) \quad \text { where } f_{j}=\left.D_{r}^{j} f\right|_{d M}=\left.(-i)^{j} \frac{\partial^{j} f}{\partial r^{j}}\right|_{d M}
$$

assigns to any smooth section its Cauchy data.
Let $W^{\prime}$ be an auxiliary graded vector bundle over $d M$. We assume that $\operatorname{dim} W=d \cdot \operatorname{dim} V$ is even and that $2 \operatorname{dim} W^{\prime}=\operatorname{dim} W$. Let $B: C^{\infty}(W) \rightarrow$ $C^{\infty}\left(W^{\prime}\right)$ be a tangential differential operator over $d M$. Decompose $B=$ $B_{i j}$ for

$$
B_{i j}: C^{\infty}\left(W_{i}\right) \rightarrow C^{\infty}\left(W_{j}^{\prime}\right)
$$

and assume that

$$
\operatorname{ord}\left(B_{i j}\right) \leq j-i
$$

It is natural to regard a section to $C^{\infty}\left(W_{i}\right)$ as being of order $i$ and to define the graded leading symbol of $B$ by:

$$
\sigma_{g}(B)_{i j}(y, \zeta)= \begin{cases}\sigma_{L}\left(B_{i j}\right)(y, \zeta) & \text { if } \operatorname{ord}\left(B_{i j}\right)=j-i \\ 0 & \text { if } \operatorname{ord}\left(B_{i j}\right)<j-i\end{cases}
$$

We then regard $\sigma_{g}(B)$ as being of graded homogeneity 0 .
We let $P_{B}$ be the operator $P$ restricted to those $f \in C^{\infty}(V)$ such that $B \underline{\gamma} f=0$. For example, let $P=-\partial^{2} / \partial x^{2}$ on the interval [ $\left.0, A\right]$. To define Dirichlet boundary conditions at $x=0$, we would set:

$$
W^{\prime}=\mathbf{C} \oplus 0 \quad \text { and } \quad B_{0,0}=1, B_{1,1}=B_{0,1}=B_{1,0}=0
$$

while to define Neumann boundary conditions at $x=0$, we would set:

$$
W^{\prime}=0 \oplus \mathbf{C} \quad \text { and } \quad B_{1,1}=1, B_{0,0}=B_{0,1}=B_{1,0}=0
$$

To define the notion of ellipticity we shall need, we consider the ordinary differential equation:

$$
p\left(y, 0, \zeta, D_{r}\right) f(r)=\lambda f(r) \quad \text { with } \lim _{r \rightarrow \infty} f(r)=0
$$

where

$$
(\zeta, \lambda) \neq(0,0) \in T^{*}(d M) \times \mathbf{C}-\mathbf{R}_{+} .
$$

We say that $(P, B)$ is elliptic with respect to $\mathbf{C}-\mathbf{R}_{+}$if $\operatorname{det}(p(x, \xi)-\lambda) \neq 0$ on the interior for all $(\zeta, \lambda) \neq(0,0) \in T^{*}(M) \times \mathbf{C}-\mathbf{R}_{+}$and if on the boundary there always exists a unique solution to this ordinary differential equation such that $\sigma_{g}(B)(y, \zeta) \underline{\gamma} f=f^{\prime}$ for any prescribed $f^{\prime} \in W^{\prime}$. In a similar fashion, we define ellipticity with respect to $\mathbf{C}-\mathbf{R}_{+}-\mathbf{R}_{-}$if these conditions hold for $\lambda \in \mathbf{C}-\mathbf{R}_{+}-\mathbf{R}_{-}$.

Again, we illustrate these notions for the operator $P=-\partial^{2} / \partial x^{2}$ and a boundary condition at $x=0$. Since $\operatorname{dim} M=1$, there is no dependence on $\zeta$ and we must simply study the ordinary differential equation:

$$
-f^{\prime \prime}=\lambda f \quad \text { with } \lim _{r \rightarrow \infty} f(r)=0 \quad(\lambda \neq 0)
$$

If $\lambda \in \mathbf{C}-\mathbf{R}_{+}$, then we can express $\lambda=\mu^{2}$ for $\operatorname{Im}(\mu)>0$. Solutions to the equation $-f^{\prime \prime}=\lambda f$ are of the form $f(r)=a e^{i \mu r}+b e^{-i \mu r}$. The decay at $\infty$ implies $b=0$ so $f(r)=a e^{i \mu r}$. Such a function is uniquely determined by either its Dirichlet or Neumann data at $r=0$ and hence $P$ is elliptic with respect to either Neumann or Dirichlet boundary conditions.

We assume henceforth that $P_{B}$ is self-adjoint and that $(P, B)$ is elliptic with respect to either the cone $\mathbf{C}-\mathbf{R}_{+}$or the cone $\mathbf{C}-\mathbf{R}_{+}-\mathbf{R}_{-}$. It is beyond the scope of this book to develop the analysis required to discuss elliptic boundary value problems; we shall simply quote the required results and refer to the appropriate papers of Seeley and Greiner for further details. Lemmas 1.6.3 and 1.6.4 generalize to yield:
Lemma 1.9.1. Let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be an elliptic partial differential operator of order $d>0$. Let $B$ be a boundary condition. We assume $(P, B)$ is self-adjoint and elliptic with respect to $\mathbf{C}-\mathbf{R}_{+}-\mathbf{R}_{-}$.
(a) We can find a complete orthonormal system $\left\{\phi_{n}\right\}_{n=1}^{\infty}$ for $L^{2}(V)$ with $P \phi_{n}=\lambda_{n} \phi_{n}$.
(b) $\phi_{n} \in C^{\infty}(V)$ and satisfy the boundary condition $B \underline{\gamma} \phi_{n}=0$.
(c) $\lambda_{n} \in \mathbf{R}$ and $\lim _{n \rightarrow \infty}\left|\lambda_{n}\right|=\infty$. If we order the $\lambda_{n}^{-}$so $\left|\lambda_{1}\right| \leq\left|\lambda_{2}\right| \cdots$ then there exists $n_{0}$ and $\delta>0$ so that $\left|\lambda_{n}\right|>n^{\delta}$ for $n>n_{0}$.
(d) If $(P, B)$ is elliptic with respect to the cone $\mathbf{C}-\mathbf{R}_{+}$, then the $\lambda_{n}$ are bounded from below and $\operatorname{spec}\left(P_{B}\right)$ is contained in $[-C, \infty)$ for some $C$.

For example, if $P=-\frac{\partial^{2}}{\partial x^{2}}$ on $[0, \pi]$ with Dirichlet boundary conditions, then the spectral resolution becomes $\left\{\sqrt{\frac{2}{\pi}} \sin (n x)\right\}_{n=1}^{\infty}$ and the corresponding eigenvalues are $n^{2}$.

If $(P, B)$ is elliptic with respect to the cone $\mathbf{C}-\mathbf{R}_{+}$, then $e^{-t P_{B}}$ is a smoothing operator with smooth kernel $K\left(t, x, x^{\prime}\right)$ defined by:

$$
K\left(t, x, x^{\prime}\right)=\sum_{n} e^{-t \lambda_{n}} \phi_{n}(x) \otimes \bar{\phi}_{n}\left(x^{\prime}\right) .
$$

Lemma 1.7.4 generalizes to yield:

Lemma 1.9.2. Let $(P, B)$ be elliptic with respect to the cone $\mathbf{C}-\mathbf{R}_{+}$, and be of order $d>0$.
(a) $e^{-t P_{B}}$ is an infinitely smoothing operator with smooth kernel $K\left(t, x, x^{\prime}\right)$.
(b) On the interior we define $a_{n}(x, P)=\operatorname{Tr} e_{n}(x, P)$ be as in Lemma 1.7.4.

On the boundary we define $a_{n}(y, P, B)$ using a complicated combinatorial recipe which depends functorially on a finite number of jets of the symbols of $P$ and of $B$.
(c) As $t \rightarrow 0^{+}$we have an asymptotic expansion:

$$
\begin{aligned}
\operatorname{Tr} e^{-t P_{B}}=\sum_{n} e^{-t \lambda_{n}}= & \int_{M} \operatorname{Tr} K(t, x, x) \operatorname{dvol}(x) \\
\sim & \sum_{n=0}^{\infty} t^{\frac{n-m}{d}} \int_{M} a_{n}(x, P) \operatorname{dvol}(x) \\
& +\sum_{n=0}^{\infty} t^{\frac{n-m+1}{d}} \int_{d M} a_{n}(y, P, B) \operatorname{dvol}(y)
\end{aligned}
$$

(d) $a_{n}(x, P)$ and $a_{n}(y, P, B)$ are invariantly defined scalar valued functions which do not depend on the coordinate system chosen nor on the local frame chosen. $a_{n}(x, P)=0$ if $n$ is odd, but $a_{n}(y, P, B)$ is in general non-zero for all values of $n$.

The interior term $a_{n}(x, P)$ arises from the calculus described previously. We first construct a parametrix on the interior. Since this parametrix will not have the proper boundary values, it is necessary to add a boundary correction term which gives rise to the additional boundary integrands $a_{n}(y, P, B)$.

To illustrate this asymptotic series, we let $P=-\frac{\partial^{2}}{\partial x^{2}}-e(x)$ on the interval $[0, A]$ where $e(x)$ is a real potential. Let $B$ be the modified Neumann boundary conditions: $f^{\prime}(0)+s(0) f(0)=f^{\prime}(A)+s(A) f(A)=0$. This is elliptic and self-adjoint with respect to the cone $\mathbf{C}-\mathbf{R}_{+}$. It can be computed that:

$$
\begin{aligned}
& \operatorname{Tr} e^{-t P_{B}} \\
& \qquad \begin{aligned}
\sim(4 \pi t)^{-1 / 2} \int\{1+t & \left.\cdot e(x)+t^{2} \cdot\left(e^{\prime \prime}(x)+3 e(x)^{2}\right) / 6+\cdots\right\} d x \\
& +\frac{1}{2}+\left(\frac{t}{\pi}\right)^{1 / 2}(s(0)-s(A)) \\
& +\frac{t}{4} \cdot\left(e(0)+e(A)+2 s^{2}(0)+2 s^{2}(A)\right)+\cdots
\end{aligned}
\end{aligned}
$$

The terms arising from the interior increase by integer powers of $t$ while the terms from the boundary increase by powers of $t^{1 / 2}$. In this integral, $d x$ is ordinary unnormalized Lebesgue measure.

In practice, we shall be interested in first order operators which are elliptic with respect to the cone $\mathbf{C}-\mathbf{R}_{+}-\mathbf{R}_{-}$and which have both positive and negative spectrum. An interesting measure of the spectral asymmetry of such an operator is obtained by studying $\operatorname{Tr}\left(P_{B} e^{-t P_{B}^{2}}\right)$ which will be discussed in more detail in the next section. To study the index of an elliptic operator, however, it suffices to study $e^{-t P_{B}^{2}}$; it is necessary to work with $P_{B}^{2}$ of course to ensure that the spectrum is positive so this converges to define a smoothing operator. There are two approaches which are available. The first is to work with the function $e^{-t \lambda^{2}}$ and integrate over a path of the form:

to define $e^{-t P_{B}^{2}}$ directly using the functional calculus. The second approach is to define $e^{-t P_{B}^{2}}$ using the operator $P^{2}$ with boundary conditions $B \underline{\gamma} f=$ $B \underline{\gamma} P f=0$. These two approaches both yield the same operator and give an appropriate asymptotic expansion which generalizes Lemma 1.9.2.

We use the heat equation to construct a local formula for the index of certain elliptic complexes. Let $Q: C^{\infty}\left(V_{1}\right) \rightarrow C^{\infty}\left(V_{2}\right)$ be an elliptic differential operator of order $d>0$. Let $Q^{*}: C^{\infty}\left(V_{2}\right) \rightarrow C^{\infty}\left(V_{1}\right)$ be the formal adjoint. Let $B_{1}: C^{\infty}\left(W_{1}\right) \rightarrow C^{\infty}\left(W_{1}^{\prime}\right)$ be a boundary condition for the operator $Q$. We assume there exists a boundary condition $B_{2}$ for $Q^{*}$ of the same form so that

$$
\left(Q_{B_{1}}\right)^{*}=\left(Q^{*}\right)_{B_{2}}
$$

We form:

$$
P=Q \oplus Q^{*} \quad \text { and } \quad B=B_{1} \oplus B_{2}
$$

so $P: C^{\infty}\left(V_{1} \oplus V_{2}\right) \rightarrow C^{\infty}\left(V_{1} \oplus V_{2}\right)$. Then $P_{B}$ will be self-adjoint. We assume that $(P, B)$ is elliptic with respect to the cone $\mathbf{C}-\mathbf{R}_{+}-\mathbf{R}_{-}$.

This is a very strong assmption which rules out many interesting cases, but is necessary to treat the index theorem using heat equation methods. We emphasize that the Atiyah-Bott theorem in its full generality does not follow from these methods.

Since $P^{2}=Q^{*} Q+Q Q^{*}$, it is clear that $P_{B}^{2}$ decomposes as the sum of two operators which preserve $C^{\infty}\left(V_{1}\right)$ and $C^{\infty}\left(V_{2}\right)$. We let $\mathcal{S}$ be the endomorphism +1 on $V_{1}$ and -1 on $V_{2}$ to take care of the signs; it is clear $P_{B}^{2} \mathcal{S}=\mathcal{S} P_{B}^{2}$. The same cancellation lemma we have used previously yields:

$$
\operatorname{index}\left(Q_{1}, B_{1}\right)=\operatorname{dim} \mathrm{N}\left(Q_{B_{1}}\right)-\operatorname{dim} \mathrm{N}\left(Q_{B_{2}}^{*}\right)=\operatorname{Tr}\left(\mathcal{S} e^{-t P_{B}^{2}}\right)
$$

(The fact that this definition of the index agrees with the definition given in the Atiyah-Bott paper follows from the fact that $B: C^{\infty}(W) \rightarrow C^{\infty}\left(W_{1}\right)$ is surjective). Consequently, the application of Lemma 1.9.2 yields:
THEOREM 1.9.3. Let $Q: C^{\infty}\left(V_{1}\right) \rightarrow C^{\infty}\left(V_{2}\right)$ be an elliptic differential operator of order $d>0$. Let $Q^{*}: C^{\infty}\left(V_{2}\right) \rightarrow C^{\infty}\left(V_{1}\right)$ be the formal adjoint and define $P=Q \oplus Q^{*}: C^{\infty}\left(V_{1} \oplus V_{2}\right) \rightarrow C^{\infty}\left(V_{1} \oplus V_{2}\right)$. Let $B=B_{1} \oplus B_{2}$ be a boundary condition such that $(P, B)$ is elliptic with respect to the cone $\mathbf{C}-\mathbf{R}_{+}-\mathbf{R}_{-}$and so that $P_{B}$ is self-adjoint. Define $\mathcal{S}=+1$ on $V_{1}$ and -1 on $V_{2}$ then:
(a) $\operatorname{index}\left(Q_{B_{1}}\right)=\operatorname{Tr}\left(\mathcal{S} e^{-t P_{B}^{2}}\right)$ for all $t$.
(b) There exist local invariants $a_{n}(x, Q)$ and $a_{n}\left(y, Q, B_{1}\right)$ such that:

$$
\begin{aligned}
\operatorname{Tr}\left(\mathcal{S} e^{-t P_{B}^{2}}\right) \sim \sum_{n=0}^{\infty} t^{\frac{n-m}{2 d}} & \int_{M} a_{n}(x, Q) \operatorname{dvol}(x) \\
& +\sum_{n=0}^{\infty} t^{\frac{n-m+1}{2 d}} \int_{d M} a_{n}\left(y, Q, B_{1}\right) \operatorname{dvol}(y)
\end{aligned}
$$

(c)

$$
\begin{aligned}
& \int_{M} a_{n}(x, Q) \operatorname{dvol}(x)+\int_{d M} a_{n-1}\left(y, Q, B_{1}\right) \operatorname{dvol}(y) \\
& \quad= \begin{cases}\operatorname{index}\left(Q_{B_{1}}\right) & \text { if } n=m \\
0 & \text { if } n \neq m\end{cases}
\end{aligned}
$$

This gives a local formula for the index; there is an analogue of Lemma 1.7.5 giving various functorial properties of these invariants we will discuss in the fourth chapter where we shall discuss the de Rham complex and the Gauss-Bonnet theorem for manifolds with boundary.

We now specialize henceforth to first order operators. We decompose

$$
p(x, \xi)=\sum_{j} e_{j}(x) \xi_{j}
$$

and near the boundary express:

$$
p(y, 0, \zeta, z)=e_{0}(y) \cdot z+\sum_{j=1}^{m-1} e_{j}(y) \zeta_{j}=e_{0}(y) \cdot z+p(y, 0, \zeta, 0)
$$

We study the ordinary differential equation:

$$
\left\{i p_{0} \partial / \partial r+p(y, 0, \zeta, 0)\right\} f(r)=\lambda f(r)
$$

or equivalently:

$$
\left.-i p_{0}\left\{\partial / \partial r+i p_{0}^{-1} p(y, 0, \zeta, 0)-i p_{0}^{-1} \lambda\right)\right\} f(r)=0
$$

With this equation in mind, we define:

$$
\tau(y, \zeta, \lambda)=i p_{0}^{-1}(p(y, 0, \zeta, 0)-\lambda)
$$

Lemma 1.9.4. Let $p(x, \xi)$ be self-adjoint and elliptic. Then $\tau(y, \zeta, \lambda)$ has no purely imaginary eigenvalues for $(\zeta, \lambda) \neq(0,0) \in T^{*}(d M) \times\left(\mathbf{C}-\mathbf{R}_{+}-\right.$ $\mathbf{R}_{-}$).
Proof: We suppose the contrary and set $\tau(y, \zeta, \lambda) v=i z v$ where $z$ is real and $v \neq 0$. This implies that:

$$
(p(y, 0, \zeta, 0)-\lambda) v=p_{0} z v
$$

or equivalently that:

$$
p(y, 0, \zeta,-z) v=\lambda v
$$

Since $p$ is self-adjoint, this implies $\lambda=0$ so $\zeta \neq 0$ which contradicts the ellipticity of $p$.

We define bundles $\Pi_{ \pm}(\tau)$ over $T^{*}(d M) \times\left\{\mathbf{C}-\mathbf{R}_{+}-\mathbf{R}_{-}\right\}-(0,0)$ to be the span of the generalized eigenvectors of $\tau$ which correspond to eigenvalues with positive/negative real part; $\Pi_{+} \oplus \Pi_{-}=V$. The differential equation has the form:

$$
\{\partial / \partial r+\tau\} f=0
$$

so the condition $\lim _{r \rightarrow \infty} f(r)=0$ implies that $f(0) \in \Pi_{+}(\tau)$. Thus $\Pi_{+}(\tau)$ is the bundle of Cauchy data corresponding to solutions to this ODE. Since $d=1$, the boundary condition $B$ is just an endomorphism:

$$
B: V_{\mid d M} \rightarrow W^{\prime}
$$

and we conclude:
Lemma 1.9.5. Let $P$ be a first order formally self-adjoint elliptic differential operator. Let $B$ be a $0^{\text {th }}$ order boundary condition. Define

$$
\tau(y, \zeta, \lambda)=i p_{0}^{-1}(p(y, 0, \zeta, 0)-\lambda) .
$$

Then $\tau$ has no purely imaginary eigenvalues and we define $\Pi_{ \pm}(\tau)$ to be bundles of generalized eigenvectors corresponding to eigenvalues with positive real and negative real parts. $(P, B)$ is elliptic with respect to the cone
$\mathbf{C}-\mathbf{R}_{+}-\mathbf{R}_{-}$if and only if $B: \Pi_{+}(\tau) \rightarrow W^{\prime}$ is an isomorphism for all $(\zeta, \lambda) \neq(0,0) \in T^{*}(d M) \times\left\{\mathbf{C}-\mathbf{R}_{+}-\mathbf{R}_{-}\right\} . \quad P_{B}$ is self-adjoint if and only if $p_{0} \mathrm{~N}(B)$ is perpendicular to $\mathrm{N}(B)$.
Proof: We have checked everything except the condition that $P_{B}$ be selfadjoint. Since $P$ is formally self-adjoint, it is immediate that:

$$
(P f, g)-(f, P g)=\int_{d M}\left(-i p_{0} f, g\right)
$$

We know that on the boundary, both $f$ and $g$ have values in $\mathrm{N}(B)$ since they satisfy the boundary condition. Thus this vanishes identically if and only if $\left(p_{0} f, g\right)=0$ for all $f, g \in \mathrm{~N}(B)$ which completes the proof.

We emphasize that these boundary conditions are much stronger than those required in the Atiyah-Bott theorem to define the index. They are much more rigid and avoid some of the pathologies which can occur otherwise. Such boundary conditions do not necessarily exist in general as we shall see later.

We shall discuss the general case in more detail in Chapter 4. We complete this section by discussing the one-dimensional case case to illustrate the ideas involved. We consider $V=[0,1] \times \mathbf{C}^{2}$ and let $P$ be the operator:

$$
P=-i \frac{\partial}{\partial r}\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right)
$$

At the point 0 , we have:

$$
\tau(\lambda)=-i \lambda\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right)
$$

Thus if $\operatorname{Im}(\lambda)>0$ we have $\Pi_{+}(\tau)=\left\{\binom{a}{a}\right\}$ while if $\operatorname{Im}(\lambda)<0$ we have $\Pi_{+}(\tau)=\left\{\binom{a}{-a}\right\}$. We let $B$ be the boundary projection which projects on the first factor- $\mathrm{N}(B)=\left\{\binom{0}{a}\right\}$; we take the same boundary condition at $x=1$ to define an elliptic self-adjoint operator $P_{B}$. Let $P^{2}=-\frac{\partial^{2}}{\partial x^{2}}$ with boundary conditions: Dirichlet boundary conditions on the first factor and Neumann boundary conditions on the second factor. The index of the problem is -1 .

### 1.10. Eta and Zeta Functions.

We have chosen to work with the heat equation for various technical reasons. However, much of the development of the subject has centered on the zeta function so we shall briefly indicate the relationship between these two in this section. We shall also define the eta invariant which plays an important role in the Atiyah-Singer index theorem for manifolds with boundary.

Recall that $\Gamma$ is defined by:

$$
\Gamma(s)=\int_{0}^{\infty} t^{s-1} e^{-t} d t
$$

for $\operatorname{Re}(s)>0$. We use the functional equation $s \Gamma(s)=\Gamma(s+1)$ to extend $\Gamma$ to a meromorphic function on $C$ with isolated simple poles at $s=0,-1,-2, \ldots$ Let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be an elliptic self-adjoint partial differential operator of order $d>0$ with positive definite leading symbol. We showed in section 1.6 that this implies $\operatorname{spec}(P)$ is contained in $[-C, \infty)$ for some constant $C$. We now assume that $P$ itself is positive definite-i.e., $\operatorname{spec}(P)$ is contained in $[\varepsilon, \infty)$ for some $\varepsilon>0$.

Proceeding formally, we define $P^{s}$ by:

$$
P^{s}=\frac{1}{2 \pi i} \int_{\alpha} \lambda^{s}(P-\lambda)^{-1} d \lambda
$$

where $\alpha$ is a suitable path in the half-plane $\operatorname{Re}(\lambda)>0$. The estimates of section 1.6 imply this is smoothing operator if $\operatorname{Re}(s) \ll 0$ with a kernel function given by:

$$
L(s, x, y)=\sum_{n} \lambda_{n}^{s} \phi_{n}(x) \otimes \bar{\phi}_{n}(y) ;
$$

this converges to define a $C^{k}$ kernel if $\operatorname{Re}(s)<s_{0}(k)$.
We use the Mellin transform to relate the zeta and heat kernels.

$$
\int_{0}^{\infty} t^{s-1} e^{-\lambda t} d t=\lambda^{-s} \int_{0}^{\infty}(\lambda t)^{s-1} e^{-\lambda t} d(\lambda t)=\lambda^{-s} \Gamma(s)
$$

This implies that:

$$
\Gamma(s) \operatorname{Tr}\left(P^{-s}\right)=\int_{0}^{\infty} t^{s-1} \operatorname{Tr} e^{-t P} d t
$$

We define $\zeta(s, P)=\operatorname{Tr}\left(P^{-s}\right)$; this is holomorphic for $\operatorname{Re}(s) \gg 0$. We decompose this integral into $\int_{0}^{1}+\int_{1}^{\infty}$. We have bounded the eigenvalues of
$P$ away from zero. Using the growth estimates of section 1.6 it is immediate that:

$$
\int_{1}^{\infty} t^{s-1} \operatorname{Tr} e^{-t P} d t=r(s, P)
$$

defines an entire function of $s$. If $0<t<1$ we use the results of section 1.7 to expand

$$
\begin{aligned}
\operatorname{Tr} e^{-t P} & =\sum_{n \leq n_{0}} t^{\frac{n-m}{d}} a_{n}(P)+O\left(t^{\frac{n_{0}-m}{d}}\right) \\
a_{n}(P) & =\int_{M} a_{n}(x, P) \operatorname{dvol}(x)
\end{aligned}
$$

where $a_{n}(x, P)$ is a local scalar invariant of the jets of the total symbol of $P$ given by Lemma 1.7.4.

If we integrate the error term which is $O\left(t^{\frac{n_{0}-m}{d}}\right)$ from 0 to 1 , we define a holomorphic function of $s$ for $\operatorname{Re}(s)+\frac{n_{0}-m}{d}>0$. We integrate $t^{s-1} t^{\frac{(n-m)}{d}}$ between 0 and 1 to conclude:

$$
\Gamma(s) \operatorname{Tr}\left(P^{-s}\right)=\Gamma(s) \zeta(s, P)=\sum_{n<n_{0}} d(s d+n-m)^{-1} a_{n}(P)+r_{n_{0}}
$$

where $r_{n_{0}}$ is holomorphic for $\operatorname{Re}(s)>-\frac{\left(n_{0}-m\right)}{d}$. This proves $\Gamma(s) \zeta(s, P)$ extends to a meromorphic function on $C$ with isolated simple poles. Furthermore, the residue at these poles is given by a local formula. Since $\Gamma$ has isolated simple poles at $s=0,-1,-2, \ldots$ we conclude that:

Lemma 1.10.1. Let $P$ be a self-adjoint, positive, elliptic partial differential operator of order $d>0$ with positive definite symbol. We define:

$$
\zeta(s, P)=\operatorname{Tr}\left(P^{-s}\right)=\Gamma(s)^{-1} \int_{0}^{\infty} t^{s-1} \operatorname{Tr} e^{-t P} d t
$$

This is well defined and holomorphic for $\operatorname{Re}(s) \gg 0$. It has a meromorphic extension to $C$ with isolated simple poles at $s=(m-n) / d$ for $n=$ $0,1,2, \ldots$. The residue of $\zeta$ at these local poles is $a_{n}(P) \Gamma((m-n) / d)^{-1}$. If $s=0,-1, \ldots$ is a non-positive integer, then $\zeta(s, P)$ is regular at this value of $s$ and its value is given by $a_{n}(P) \operatorname{Res}_{s=(m-n) / d} \Gamma(s) . a_{n}(P)$ is the invariant given in the asymptotic expansion of the heat equation. $a_{n}(P)=\int_{M} a_{n}(x, P) \mathrm{dvol}(x)$ where $a_{n}(x, P)$ is a local invariant of the jets of the total symbol of $P . a_{n}$ vanishes if $n$ is odd.

Remark: If $A$ is an auxilary differential operator of order $a$, we can define

$$
\zeta(s, A, P)=\operatorname{Tr}\left(A P^{-s}\right)=\Gamma(s)^{-1} \int_{0}^{\infty} t^{s-1} \operatorname{Tr}\left(A e^{-t P}\right) d t
$$

We apply Lemma 1.7.7 to see this has a meromorphic extension to $\mathbf{C}$ with isolated simple poles at $s=(m+a-n) / d$. The residue at these poles is given by the generalized invariants of the heat equation

$$
\frac{a_{n}(A, P)}{\Gamma\{(m+a-n) / d\}} .
$$

There is a similar theorem if $d M \neq \emptyset$ and if $B$ is an elliptic boundary condition as discussed in section 1.9.

If $P$ is only positive semi-definite so it has a finite dimensional space corresponding to the eigenvalue 0 , we define:

$$
\zeta(s, P)=\sum_{\lambda_{n}>0} \lambda_{n}^{-s}
$$

and Lemma 1.10 extends to this case as well with suitable modifications. For example, if $M=S^{1}$ is the unit circle and if $P=-\partial^{2} / \partial \theta^{2}$ is the Laplacian, then:

$$
\zeta(s, P)=2 \sum_{n>0} n^{-2 s}
$$

is essentially just the Riemann zeta function. This has a simple isolated pole at $s=1 / 2$.

If $Q: C^{\infty}\left(V_{1}\right) \rightarrow C^{\infty}\left(V_{2}\right)$, then the cancellation lemma of section 1.6 implies that:

$$
\varepsilon^{-s} \operatorname{index}(Q)=\zeta\left(s, Q^{*} Q+\varepsilon\right)-\zeta\left(s, Q Q^{*}+\varepsilon\right)
$$

for any $\varepsilon>0$ and for any $s . \zeta$ is regular at $s=0$ and its value is given by a local formula. This gives a local formula for index $(Q)$. Using Lemma 1.10.1 and some functorial properties of the heat equation, it is not difficult to show this local formula is equivalent to the local formula given by the heat equation so no new information has resulted. The asymptotics of the heat equation are related to the values and residues of the zeta function by Lemma 1.10.1.

If we do not assume that $P$ is positive definite, it is possible to define a more subtle invariant which measures the difference between positive and negative spectrum. Let $P$ be a self-adjoint elliptic partial differential operator of order $d>0$ which is not necessarily positive definite. We define the eta invariant by:

$$
\begin{aligned}
\eta(s, P) & =\sum_{\lambda_{n}>0} \lambda_{n}^{-s}-\sum_{\lambda_{n}<0}\left(-\lambda_{n}\right)^{-s}=\sum_{\lambda_{n} \neq 0} \operatorname{sign}\left(\lambda_{n}\right)\left|\lambda_{n}\right|^{-s} \\
& =\operatorname{Tr}\left(P \cdot\left(P^{2}\right)^{-(s+1) / 2}\right) .
\end{aligned}
$$

Again, this is absolutely convergent and defines a holomorphic function if $\operatorname{Re}(s) \gg 0$.

We can also discuss the eta invariant using the heat equation. The identity:

$$
\int_{0}^{\infty} t^{(s-1) / 2} \lambda e^{-\lambda^{2} t} d t=\Gamma((s+1) / 2) \operatorname{sign}(\lambda)|\lambda|^{-s}
$$

implies that:

$$
\eta(s, P)=\Gamma((s+1) / 2)^{-1} \int_{0}^{\infty} t^{(s-1) / 2} \operatorname{Tr}\left\{P e^{-t P^{2}}\right\} d t
$$

Again, the asymptotic behavior at $t=0$ is all that counts in producing poles since this decays exponentially at $\infty$ assuming $P$ has no zero eigenvector; if $\operatorname{dim} \mathrm{N}(P)>0$ a seperate argument must be made to take care of this eigenspace. This can be done by replacing $P$ by $P+\varepsilon$ and letting $\varepsilon \rightarrow 0$.

Lemma 1.7.7 shows that there is an asymptotic series of the form:

$$
\operatorname{Tr}\left(P e^{-t P^{2}}\right) \sim \sum_{n=0}^{\infty} t^{\frac{n-m-d}{2 d}} a_{n}\left(P, P^{2}\right)
$$

for

$$
a_{n}\left(P, P^{2}\right)=\int_{M} a_{n}\left(x, P, P^{2}\right) \operatorname{dvol}(x)
$$

This is a local invariant of the jets of the total symbol of $P$.
We substitute this asymptotic expansion into the expression for $\eta$ to see:

$$
\eta(s, P) \Gamma((s+1) / 2)^{-1}=\sum_{n \leq n_{0}} \frac{2 d}{d s+n-m} a_{n}\left(P, P^{2}\right)+r_{n_{0}}
$$

where $r_{n_{0}}$ is holomorphic on a suitable half-plane. This proves $\eta$ has a suitable meromorphic extension to $\mathbf{C}$ with locally computable residues.

Unfortunately, while it was clear from the analysis that $\zeta$ was regular at $s=0$, it is not immediate that $\eta$ is regular at $s=0$ since $\Gamma$ does not have a pole at $\frac{1}{2}$ to cancel the pole which may be introduced when $n=m$. in fact, if one works with the local invariants involved, $a_{n}\left(x, P, P^{2}\right) \neq 0$ in general so the local poles are in fact present at $s=0$. However, it is a fact which we shall discuss later that $\eta$ is regular at $s=0$ and we will define

$$
\tilde{\eta}(p)=\frac{1}{2}\{\eta(0, p)+\operatorname{dim} \mathrm{N}(p)\} \quad \bmod \mathbf{Z}
$$

(we reduce modulo $\mathbf{Z}$ since $\eta$ has jumps in $2 \mathbf{Z}$ as eigenvalues cross the origin).

We compute a specific example to illustrate the role of $\eta$ in measuring spectral asymmetry. Let $P=-i \partial / \partial \theta$ on $C^{\infty}\left(S^{1}\right)$, then the eigenvalues of $p$ are the integers so $\eta(s, p)=0$ since the spectrum is symmetric about the origin. Let $a \in \mathbf{R}$ and define:

$$
P_{a}=P-a, \quad \eta\left(s, P_{a}\right)=\sum_{n \in \mathbf{Z}} \operatorname{sign}\{n-a\}|n-a|^{-s}
$$

We differentiate this with respect to the parameter $a$ to conclude:

$$
\frac{d}{d a} \eta(s, P(a))=\sum_{n \in \mathbf{Z}} s\left((n-a)^{2}\right)^{-(s+1) / 2}
$$

If we compare this with the Riemann zeta function, then

$$
\sum_{n \in \mathbf{Z}}\left((n-a)^{2}\right)^{-(s+1) / 2}
$$

has a simple pole at $s=0$ with residue 2 and therefore:

$$
\left.\frac{d}{d a} \eta(s, P(a))\right|_{s=0}=2
$$

Since $\eta$ vanishes when $a=0$, we integrate this with respect to $a$ to conclude:

$$
\eta(0, P(a))=2 a+1 \bmod 2 \mathbf{Z} \quad \text { and } \quad \tilde{\eta}(P(a))=a+\frac{1}{2}
$$

is non-trivial. (We must work modulo $\mathbf{Z}$ since spec $P(a)$ is periodic with period 1 in $a)$.

We used the identity:

$$
\frac{d}{d a} \eta\left(s, P_{a}\right)=-s \operatorname{Tr}\left(\dot{P}_{a}\left(P_{a}^{2}\right)^{-(s+1) / 2}\right)
$$

in the previous computation; it in fact holds true in general:
LEmMA 1.10.2. Let $P(a)$ be a smooth 1-parameter family of elliptic selfadjoint partial differential operators of order $d>0$. Assume $P(a)$ has no zero spectrum for $a$ in the parameter range. Then if "." denotes differentiation with respect to the parameter $a$,

$$
\dot{\eta}(s, P(a))=-s \operatorname{Tr}\left(\dot{P}(a)\left(P(a)^{2}\right)^{-(s+1) / 2}\right)
$$

If $P(a)$ has zero spectrum, we regard $\eta(s, P(a)) \in \mathbf{C} / \mathbf{Z}$.
Proof: If we replace $P$ by $P^{k}$ for an odd positive integer $k$, then $\eta\left(s, P^{k}\right)$ $=\eta(k s, P)$. This shows that it suffices to prove Lemma 1.10.2 for $d \gg 0$.

By Lemma 1.6.6, $(P-\lambda)^{-1}$ is smoothing and hence of trace class for $d$ large. We can take trace inside the integral to compute:

$$
\begin{aligned}
& \eta(s, P(a))=\frac{1}{2 \pi i}\left\{\int_{\gamma_{1}} \lambda^{-s} \operatorname{Tr}\left((P(a)-\lambda)^{-1}\right) d \lambda\right. \\
&-\int_{\gamma_{2}}(-\lambda)^{-s} \operatorname{Tr}\left((P(a)-\lambda)^{-1} d \lambda\right\}
\end{aligned}
$$

where $\gamma_{1}$ and $\gamma_{2}$ are paths about the positive/negative real axis which are oriented suitably. We differentiate with respect to the parameter $a$ to express:

$$
\frac{d}{d a}\left((P(a)-\lambda)^{-1}=-(P(a)-\lambda)^{-1} \dot{P}(a)(P(a)-\lambda)^{-1} .\right.
$$

Since the operators involved are of Trace class,

$$
\operatorname{Tr}\left(\frac{d}{d a}(P(a)-\lambda)^{-1}\right)=-\operatorname{Tr}\left(\dot{P}(a)(P(a)-\lambda)^{-2}\right)
$$

We use this identity and bring trace back outside the integral to compute:

$$
\begin{aligned}
\eta(s, P(a))=\frac{1}{2 \pi i} \operatorname{Tr}(\dot{P}(a))\left\{\int_{\gamma_{1}}\right. & -\lambda^{-s}(P(a)-\lambda)^{-2} d \lambda \\
& \left.+\int_{\gamma_{2}}(-\lambda)^{-s}(P(a)-\lambda)^{-2} d \lambda\right\}
\end{aligned}
$$

We now use the identity:

$$
\frac{d}{d \lambda}\left((P(a)-\lambda)^{-1}\right)=(P(a)-\lambda)^{-2}
$$

to integrate by parts in $\lambda$ in this expression. This leads immediately to the desired formula.

We could also calculate using the heat equation. We proceed formally:

$$
\begin{aligned}
\frac{d \eta}{d a}(s, P(a)) & =\Gamma\{(s+1) / 2\}^{-1} \int_{0}^{\infty} t^{(s-1) / 2} \operatorname{Tr}\left(\frac{d}{d a}\left(P(a) e^{-t P(a)^{2}}\right)\right) d t \\
& =\Gamma\{(s+1) / 2\}^{-1} \int_{0}^{\infty} t^{(s-1) / 2} \operatorname{Tr}\left(\frac{d P}{d a}\left(1-2 t P^{2}\right) e^{-t P^{2}}\right) d t \\
& =\Gamma\{(s+1) / 2\}^{-1} \int_{0}^{\infty} t^{(s-1) / 2}\left(1+2 t \frac{d}{d t}\right) \operatorname{Tr}\left(\frac{d P}{d a} e^{-t P^{2}}\right) d t
\end{aligned}
$$

We now integrate by parts to compute:

$$
\begin{aligned}
& =\Gamma\{(s+1) / 2\}^{-1} \int_{0}^{\infty}-s t^{(s-1) / 2} \operatorname{Tr}\left(\frac{d P}{d a} e^{-t P^{2}}\right) d t \\
& \sim-s \Gamma\{(s+1) / 2\}^{-1} \sum_{n} \int_{0}^{1} t^{(s-1) / 2} t^{(n-m-d) / 2 d} a_{n}\left(\frac{d P}{d a}, P^{2}\right) d t \\
& \sim-s \Gamma\{(s+1) / 2\}^{-1} \sum_{n} 2 /(s+(n-m) / d) a_{n}\left(\frac{d P}{d a}, P^{2}\right)
\end{aligned}
$$

In particular, this shows that $\frac{d \tilde{\eta}}{d a}$ is regular at $s=0$ and the value

$$
-\Gamma\left(\frac{1}{2}\right)^{-1} a_{m}\left(\frac{d P}{d a}, P^{2}\right)
$$

is given by a local formula.
The interchange of order involved in using global trace is an essential part of this argument. $\operatorname{Tr}\left(\dot{P}(a)\left(P^{2}\right)^{-(s+1) / 2}\right)$ has a meromorphic extension to $\mathbf{C}$ with isolated simple poles. The pole at $s=-1 / 2$ can be present, but is cancelled off by the factor of $-s$ which multiplies the expression. Thus $\frac{d}{d a} \operatorname{Res}_{s=0} \eta(s, P(a))=0$. This shows the global residue is a homotopy invariant; this fact will be used in Chapter 4 to show that in fact $\eta$ is regular at the origin. This argument does not go through locally, and in fact it is possible to construct operators in any dimension $m \geq 2$ so that the local eta function is not regular at $s=0$.

We have assumed that $P(a)$ has no zero spectrum. If we supress a finite number of eigenvalues which may cross the origin, this makes no contribution to the residue at $s=0$. Furthermore, the value at $s=0$ changes by jumps of an even integer as eigenvalues cross the origin. This shows $\tilde{\eta}$ is in fact smooth in the parameter $a$ and proves:

## Lemma 1.10.3.

(a) Let $P$ be a self-adjoint elliptic partial differential operator of positive order d. Define:
$\eta(s, P)=\operatorname{Tr}\left(P\left(P^{2}\right)^{-(s+1) / 2}\right)=\Gamma\{(s+1) / 2\}^{-1} \int_{0}^{\infty} t^{(s-1) / 2} \operatorname{Tr}\left(P e^{-t P^{2}}\right) d t$.
This admits a meromorphic extension to $\mathbf{C}$ with isolated simple poles at $s=(m-n) / d$. The residue of $\eta$ at such a pole is computed by:

$$
\operatorname{Res}_{s=(m-n) / d}=2 \Gamma\{(m+d-n) / 2 d\}^{-1} \int_{M} a_{n}\left(x, P, P^{2}\right) \operatorname{dvol}(x)
$$

where $a_{n}\left(x, P, P^{2}\right)$ is defined in Lemma 1.7.7; it is a local invariant in the jets of the symbol of $P$.
(b) Let $P(a)$ be a smooth 1-parameter family of such operators. If we assume eigenvalues do not cross the origin, then:

$$
\begin{aligned}
\frac{d}{d a} \eta(s, P(a)) & =-s \operatorname{Tr}\left(\frac{d P}{d a} \cdot\left(P(a)^{2}\right)^{-(s+1) / 2}\right) \\
& =-2 s \Gamma\{(s+1) / 2\}^{-1} \int_{0}^{\infty} t^{(s-1) / 2} \operatorname{Tr}\left(\frac{d P}{d a} \cdot e^{-t P^{2}}\right) d t
\end{aligned}
$$

Regardless of whether or not eigenvalues cross the origin, $\tilde{\eta}(P(a))$ is smooth in $\mathbf{R} \bmod \mathbf{Z}$ in the parameter $a$ and

$$
\begin{aligned}
\frac{d}{d a} \tilde{\eta}(P(a)) & =-\Gamma\left(\frac{1}{2}\right)^{-1} a_{m}\left(\frac{d P}{d a}, P^{2}\right) \\
& =-\Gamma\left(\frac{1}{2}\right)^{-1} \int_{M} a_{m}\left(x, \frac{d P}{d a}, P^{2}\right) \operatorname{dvol}(x)
\end{aligned}
$$

is the local invariant in the jets of the symbols of $\left(\frac{d P}{d a}, P^{2}\right)$ given by Lemma 1.7.7.

In the example on the circle, the operator $P(a)$ is locally isomorphic to $P$. Thus the value of $\eta$ at the origin is not given by a local formula. This is a global invariant of the operator; only the derivative is locally given.

It is not necessary to assume that $P$ is a differential operator to define the eta invariant. If $P$ is an elliptic self-adjoint pseudo-differential operator of order $d>0$, then the sum defining $\eta(s, P)$ converges absolutely for $s \gg 0$ to define a holomorphic function. This admits a meromorphic extension to the half-plane $\operatorname{Re}(s)>-\delta$ for some $\delta>0$ and the results of Lemma 1.10.3 continue to apply. This requires much more delicate estimates than we have developed and we shall omit details. The reader is referred to the papers of Seeley for the proofs.

We also remark that it is not necessary to assume that $P$ is self-adjoint; it suffices to assume that $\operatorname{det}(p(x, \xi)-i t) \neq 0$ for $(\xi, t) \neq(0,0) \in T^{*} M \times \mathbf{R}$. Under this ellipticity hypothesis, the spectrum of $P$ is discrete and only a finite number of eigenvalues lie on or near the imaginary axis. We define:

$$
\begin{aligned}
\eta(s, P) & =\sum_{\operatorname{Re}\left(\lambda_{i}\right)>0}\left(\lambda_{i}\right)^{-s}-\sum_{\operatorname{Re}\left(\lambda_{i}\right)<0}\left(-\lambda_{i}\right)^{-s} \\
\tilde{\eta}(P) & =\frac{1}{2}\left\{\eta(s, P)-\frac{1}{s} \operatorname{Res}_{s=0} \eta(s, P)+\operatorname{dim} \mathrm{N}(i \mathbf{R})\right\}_{s=0} \quad \bmod \mathbf{Z}
\end{aligned}
$$

where $\operatorname{dim} \mathrm{N}(i \mathbf{R})$ is the dimension of the finite dimensional vector space of generalized eigenvectors of $P$ corresponding to purely imaginary eigenvalues.

In section 4.3, we will discuss the eta invariant in further detail and use it to define an index with coefficients in a locally flat bundle using secondary characteristic classes.

If the leading symbol of $P$ is positive definite, the asymptotics of $\operatorname{Tr}\left(e^{-t P}\right)$ as $t \rightarrow 0^{+}$are given by local formulas integrated over the manifold. Let $A(x)=a_{0} x^{a}+\cdots+a_{a}$ and $B(x)=b_{0} x^{b}+\cdots+b_{b}$ be polynomials where $b_{0}>$ 0 . The operator $A(P) e^{-t B(P)}$ is infinitely smoothing. The asymptotics of $\operatorname{Tr}\left(A(P) e^{-t B(P)}\right)$ are linear combinations of the invariants $a_{N}(P)$ giving the asymptotics of $\operatorname{Tr}\left(e^{-t P}\right)$. Thus there is no new information which results by considering more general operators of heat equation type. If the leading symbol of $P$ is indefinite, one must consider both the zeta and eta function or equivalently:

$$
\operatorname{Tr}\left(e^{-t P^{2}}\right) \quad \text { and } \quad \operatorname{Tr}\left(P e^{-t P^{2}}\right)
$$

Then if $b_{0}>0$ is even, we obtain enough invariants to calculate the asymptotics of $\operatorname{Tr}\left(A(P) e^{-t B(P)}\right)$. We refer to Fegan-Gilkey for further details.

## CHAPTER 2

## CHARACTERISTIC CLASSES

## Introduction

In the second chapter, we develop the theory of characteristic classes. In section 2.1, we discuss the Chern classes of a complex vector bundle and in section 2.2 we discuss the Pontrjagin and Euler classes of a real vector bundle. We shall define the Todd class, the Hirzebruch $L$-polynomial, and the $A$-roof genus which will play a central role in our discussion of the index theorem. We also discuss the total Chern and Pontrjagin classes as well as the Chern character.

In section 2.3, we apply these ideas to the tangent space of a real manifold and to the holomorphic tangent space of a holomorphic manifold. We compute several examples defined by Clifford matrices and compute the Chern classes of complex projective space. We show that suitable products of projective spaces form a dual basis to the space of characteristic classes. Such products will be used in chapter three to find the normalizing constants which appear in the formula for the index theorem.

In section 2.4 and in the first part of section 2.5 , we give a heat equation proof of the Gauss-Bonnet theorem. This is based on first giving an abstract characterization of the Euler form in terms of invariance theory. This permis us to identify the integrand of the heat equation with the Euler integrand. This gives a more direct proof of Theorem 2.4.8 which was first proved by Patodi using a complicated cancellation lemma. (The theorem for dimension $m=2$ is due to McKean and Singer).

In the remainder of section 2.5, we develop a similar characterization of the Pontrjagin forms of the real tangent space. We shall wait until the third chapter to apply these results to obtain the Hirzebruch signature theorem. There are two different approaches to this result. We have presented both our original approach and also one modeled on an approach due to Atiyah, Patodi and Bott. This approach uses H. Weyl's theorem on the invariants of the orthogonal group and is not self-contained as it also uses facts about geodesic normal coordinates we shall not develop. The other approach is more combinatorial in flavor but is self-contained. It also generalizes to deal with the holomorphic case for a Kaehler metric.

The signature complex with coefficients in a bundle $V$ gives rise to invariants which depend upon both the metric on the tangent space of $M$ and on the connection 1 -form of $V$. In section 2.6 , we extend the results of section 2.5 to cover more general invariants of this type. We also construct dual bases for these more general invariants similar to those constructed in section 2.3 using products of projective spaces and suitable bundles over these spaces.

The material of sections 2.1 through 2.3 is standard and reviews the theory of characteristic classes in the context we shall need. Sections 2.4 through 2.6 deal with less standard material. The chapter is entirely selfcontained with the exception of some material in section 2.5 as noted above. We have postponed a discussion of similar material for the holomorphic Kaehler case until sections 3.6 and 3.7 of chapter three since this material is not needed to discuss the signature and spin complexes.

### 2.1. Characteristic Classes <br> Of a Complex Vector Bundle.

The characteristic classes are topological invariants of a vector bundle which are represented by differential forms. They are defined in terms of the curvature of a connection.

Let $M$ be a smooth compact manifold and let $V$ be a smooth complex vector bundle of dimension $k$ over $M$. A connection $\nabla$ on $V$ is a first order partial differential operator $\nabla: C^{\infty}(V) \rightarrow C^{\infty}\left(T^{*} M \otimes V\right)$ such that:

$$
\nabla(f s)=d f \otimes s+f \nabla s \quad \text { for } f \in C^{\infty}(M) \text { and } s \in C^{\infty}(V)
$$

Let $\left(s_{1}, \ldots, s_{k}\right)$ be a local frame for $V$. We can decompose any section $s \in C^{\infty}(V)$ locally in the form $s(x)=f_{i}(x) s_{i}(x)$ for $f_{i} \in C^{\infty}(M)$. We adopt the convention of summing over repeated indices unless otherwise specified in this subsection. We compute:

$$
\nabla s=d f_{i} \otimes s_{i}+f_{i} \nabla s_{i}=d f_{i} \otimes s_{i}+f_{i} \omega_{i j} \otimes s_{j} \quad \text { where } \nabla s_{i}=\omega_{i j} \otimes s_{j}
$$

The connection 1-form $\omega$ defined by

$$
\omega=\omega_{i j}
$$

is a matrix of 1 -forms. The connection $\nabla$ is uniquely determined by the derivation property and by the connection 1-form. If we specify $\omega$ arbitrarily locally, we can define $\nabla$ locally. The convex combination of connections is again a connection, so using a partition of unity we can always construct connections on a bundle $V$.

If we choose another frame for $V$, we can express $s_{i}^{\prime}=h_{i j} s_{j}$. If $h_{i j}^{-1} s_{j}^{\prime}=$ $s_{i}$ represents the inverse matrix, then we compute:

$$
\begin{aligned}
\nabla s_{i}^{\prime} & =\omega_{i j}^{\prime} \otimes s_{j}^{\prime}=\nabla\left(h_{i k} s_{k}\right)=d h_{i k} \otimes s_{k}+h_{i k} \omega_{k l} \otimes s_{l} \\
& =\left(d h_{i k} h_{k j}^{-1}+h_{i k} \omega_{k l} h_{l j}^{-1}\right) \otimes s_{j}^{\prime} .
\end{aligned}
$$

This shows $\omega^{\prime}$ obeys the transformation law:

$$
\omega_{i j}^{\prime}=d h_{i k} h_{k j}^{-1}+h_{i k} \omega_{k l} h_{l j}^{-1} \quad \text { i.e., } \omega^{\prime}=d h \cdot h^{-1}+h \omega h^{-1}
$$

in matrix notation. This is, of course, the manner in which the $0^{\text {th }}$ order symbol of a first order operator transforms.

We extend $\nabla$ to be a derivation mapping

$$
C^{\infty}\left(\Lambda^{p} T^{*} M \otimes V\right) \rightarrow C^{\infty}\left(\Lambda^{p+1} T^{*} M \otimes V\right)
$$

so that:

$$
\nabla\left(\theta_{p} \otimes s\right)=d \theta_{p} \otimes s+(-1)^{p} \theta_{p} \wedge \nabla s
$$

We compute that:
$\nabla^{2}(f s)=\nabla(d f \otimes s+f \nabla s)=d^{2} f \otimes s-d f \wedge \nabla s+d f \wedge \nabla s+f \nabla^{2} s=f \nabla^{2} s$
so instead of being a second order operator, $\nabla^{2}$ is a $0^{\text {th }}$ order operator. We may therefore express

$$
\nabla^{2}(s)\left(x_{0}\right)=\Omega\left(x_{0}\right) s\left(x_{0}\right)
$$

where the curvature $\Omega$ is a section to the bundle $\Lambda^{2}\left(T^{*} M\right) \otimes \operatorname{END}(V)$ is a 2 -form valued endomorphism of $V$.

In local coordinates, we compute:

$$
\Omega\left(s_{i}\right)=\Omega_{i j} \otimes s_{j}=\nabla\left(\omega_{i j} \otimes s_{j}\right)=d \omega_{i j} \otimes s_{j}-\omega_{i j} \wedge \omega_{j k} \otimes s_{k}
$$

so that:

$$
\Omega_{i j}=d \omega_{i j}-\omega_{i k} \wedge \omega_{k j} \quad \text { i.e., } \Omega=d \omega-\omega \wedge \omega \text {. }
$$

Since $\nabla^{2}$ is a $0^{\text {th }}$ order operator, $\Omega$ must transform like a tensor:

$$
\Omega_{i j}^{\prime}=h_{i k} \Omega_{k l} h_{l j}^{-1} \quad \text { i.e., } \Omega^{\prime}=h \Omega h^{-1}
$$

This can also be verified directly from the transition law for $\omega^{\prime}$. The reader should note that in some references, the curvature is defined by $\Omega=d \omega+$ $\omega \wedge \omega$. This sign convention results from writing $V \otimes T^{*} M$ instead of $T^{*} M \otimes V$ and corresponds to studying left invariant rather than right invariant vector fields on $\mathrm{GL}(k, C)$.

It is often convenient to normalize the choice of frame.
Lemma 2.1.1. Let $\nabla$ be a connection on a vector bundle $V$. We can always choose a frame $s$ so that at a given point $x_{0}$ we have

$$
\omega\left(x_{0}\right)=0 \quad \text { and } \quad d \Omega\left(x_{0}\right)=0
$$

Proof: We find a matrix $h(x)$ defined near $x_{0}$ so that $h\left(x_{0}\right)=I$ and $d h\left(x_{0}\right)=-\omega\left(x_{0}\right)$. If $s_{i}^{\prime}=h_{i j} s_{j}$, then it is immediate that $\omega^{\prime}\left(x_{0}\right)=0$. Similarly we compute $d \Omega^{\prime}\left(x_{0}\right)=d\left(d \omega^{\prime}-\omega^{\prime} \wedge \omega^{\prime}\right)\left(x_{0}\right)=\omega^{\prime}\left(x_{0}\right) \wedge d \omega^{\prime}\left(x_{0}\right)-$ $d \omega^{\prime}\left(x_{0}\right) \wedge \omega^{\prime}\left(x_{0}\right)=0$.

We note that as the curvature is invariantly defined, we cannot in general find a parallel frame $s$ so $\omega$ vanishes in a neighborhood of $x_{0}$ since this would imply $\nabla^{2}=0$ near $x_{0}$ which need not be true.

Let $A_{i j}$ denote the components of a matrix in $\operatorname{END}\left(\mathbf{C}^{k}\right)$ and let $P(A)=$ $P\left(A_{i j}\right)$ be a polynomial mapping $\operatorname{END}\left(\mathbf{C}^{k}\right) \rightarrow \mathbf{C}$. We assume that $P$ is invariant-i.e., $P\left(h A h^{-1}\right)=P(A)$ for any $h \in \mathrm{GL}(k, \mathbf{C})$. We define $P(\Omega)$ as an even differential form on $M$ by substitution. Since $P$ is invariant
and since the curvature transforms tensorially, $P(\Omega) \stackrel{\text { def }}{=} P(\nabla)$ is defined independently of the particular local frame which is chosen.

There are two examples which are of particular interest and which will play an important role in the statement of the Atiyah-Singer index theorem. We define:

$$
\begin{aligned}
c(A) & =\operatorname{det}\left(I+\frac{i}{2 \pi} A\right) & & \\
& =1+c_{1}(A)+\cdots+c_{k}(A) & & \text { (the total Chern form) } \\
c h(A) & =\operatorname{Tr} e^{i A / 2 \pi}=\sum_{j} \operatorname{Tr}\left(\frac{i A}{2 \pi}\right)^{j} / j! & & \text { (the total Chern character) }
\end{aligned}
$$

The $c_{j}(A)$ represent the portion of $c(A)$ which is homogeneous of order $j$. $c_{j}(\nabla) \in \Lambda^{2 j}(M)$. In a similar fashion, we decompose $c h(A)=\sum_{j} c h_{j}(A)$ for $c h_{j}(A)=\operatorname{Tr}(i A / 2 \pi)^{j} / j$ !

Strictly speaking, $\operatorname{ch}(A)$ is not a polynomial. However, when we substitute the components of the curvature tensor, this becomes a finite sum since $\operatorname{Tr}\left(\Omega^{j}\right)=0$ if $2 j>\operatorname{dim} M$. More generally, we can define $P(\Omega)$ if $P(A)$ is an invariant formal power series by truncating the power series appropriately.

As a differential form, $P(\nabla)$ depends on the connection chosen. We show $P(\nabla)$ is always closed. As an element of de Rham cohomology, $P(\nabla)$ is independent of the connection and defines a cohomology class we shall denote by $P(V)$.

Lemma 2.1.2. Let $P$ be an invariant polynomial.
(a) $d P(\nabla)=0$ so $P(\nabla)$ is a closed differential form.
(b) Given two connections $\nabla_{0}$ and $\nabla_{1}$, we can define a differential form $T P\left(\nabla_{0}, \nabla_{1}\right)$ so that $P\left(\nabla_{1}\right)-P\left(\nabla_{0}\right)=d\left\{T P\left(\nabla_{1}, \nabla_{0}\right)\right\}$.

Proof: By decomposing $P$ as a sum of homogeneous polynomials, we may assume without loss of generality that $P$ is homogeneous of order $k$. Let $P\left(A_{1}, \ldots, A_{k}\right)$ denote the complete polarization of $P$. We expand $P\left(t_{1} A_{1}+\right.$ $\left.\cdots+t_{k} A_{k}\right)$ as a polynomial in the variables $\left\{t_{j}\right\} .1 / k!$ times the coefficient of $t_{1} \ldots t_{k}$ is the polarization of $P . P$ is a multi-linear symmetric function of its arguments. For example, if $P(A)=\operatorname{Tr}\left(A^{3}\right)$, then the polarization is given by $\frac{1}{2} \operatorname{Tr}\left(A_{1} A_{2} A_{3}+A_{2} A_{1} A_{3}\right)$ and $P(A)=P(A, A, A)$.

Fix a point $x_{0}$ of $M$ and choose a frame so $\omega\left(x_{0}\right)=d \Omega\left(x_{0}\right)=0$. Then

$$
d P(\Omega)\left(x_{0}\right)=d P(\Omega, \ldots, \Omega)\left(x_{0}\right)=k P(d \Omega, \Omega, \ldots, \Omega)\left(x_{0}\right)=0
$$

Since $x_{0}$ is arbitrary and since $d P(\Omega)$ is independent of the frame chosen this proves $d P(\Omega)=0$ which proves (a).

Let $\nabla_{t}=t \nabla_{1}+(1-t) \nabla_{0}$ have connection 1-form $\omega_{t}=\omega_{0}+t \theta$ where $\theta=\omega_{1}-\omega_{0}$. The transformation law for $\omega$ implies relative to a new frame:

$$
\begin{aligned}
\theta^{\prime}=\omega_{1}^{\prime}-\omega_{0}^{\prime} & =\left(d h \cdot h^{-1}+h \omega_{1} h^{-1}\right)-\left(d h \cdot h^{-1}+h \omega_{0} h^{-1}\right) \\
& =h\left(\omega_{1}-\omega_{0}\right) h^{-1}=h \theta h^{-1}
\end{aligned}
$$

so $\theta$ transforms like a tensor. This is of course nothing but the fact that the difference between two first order operators with the same leading symbol is a $0^{\text {th }}$ order operator.

Let $\Omega_{t}$ be the curvature of the connection $\nabla_{t}$. This is a matrix valued 2 -form. Since $\theta$ is a 1 -form, it commutes with 2 -forms and we can define

$$
P\left(\theta, \Omega_{t}, \ldots, \Omega_{t}\right) \in \Lambda^{2 k-1}\left(T^{*} M\right)
$$

by substitution. Since $P$ is invariant, its complete polarization is also invariant so $P\left(h^{-1} \theta h, h^{-1} \Omega_{t} h, \ldots, h^{-1} \Omega_{t} h\right)=P\left(\theta, \Omega_{t}, \ldots, \Omega_{t}\right)$ is invariantly defined independent of the choice of frame.

We compute that:

$$
P\left(\nabla_{1}\right)-P\left(\nabla_{0}\right)=\int_{0}^{1} \frac{d}{d t} P\left(\Omega_{t}, \ldots, \Omega_{t}\right) d t=k \int_{0}^{1} P\left(\Omega_{t}^{\prime}, \Omega_{t}, \ldots, \Omega_{t}\right) d t
$$

We define:

$$
T P\left(\nabla_{1}, \nabla_{0}\right)=k \int_{0}^{1} P\left(\theta, \Omega_{t}, \ldots, \Omega_{t}\right) d t
$$

To complete the proof of the Lemma, it suffices to check

$$
d P\left(\theta, \Omega_{t}, \ldots, \Omega_{t}\right)=P\left(\Omega_{t}^{\prime}, \Omega_{t}, \ldots, \Omega_{t}\right)
$$

Since both sides of the equation are invariantly defined, we can choose a suitable local frame to simplify the computation. Let $x_{0} \in M$ and fix $t_{0}$. We choose a frame so $\omega_{t}\left(x_{0}, t_{0}\right)=0$ and $d \Omega_{t}\left(x_{0}, t_{0}\right)=0$. We compute:

$$
\begin{aligned}
\Omega_{t}^{\prime} & =\left\{d \omega_{0}+t d \theta-\omega_{t} \wedge \omega_{t}\right\}^{\prime} \\
& =d \theta-\omega_{t}^{\prime} \wedge \omega_{t}-\omega_{t} \wedge \omega_{t}^{\prime} \\
\Omega_{t}^{\prime}\left(x_{0}, t_{0}\right) & =d \theta
\end{aligned}
$$

and

$$
\begin{aligned}
d P\left(\theta, \Omega_{t}, \ldots, \Omega_{t}\right)\left(x_{0}, t_{0}\right) & =P\left(d \theta, \Omega_{t}, \ldots, \Omega_{t}\right)\left(x_{0}, t_{0}\right) \\
& =P\left(\Omega_{t}^{\prime}, \Omega_{t}, \ldots, \Omega_{t}\right)\left(x_{0}, t_{0}\right)
\end{aligned}
$$

which completes the proof.
$T P$ is called the transgression of $P$ and will play an important role in our discussion of secondary characteristic classes in Chapter 4 when we discuss the eta invariant with coefficients in a locally flat bundle.

We suppose that the matrix $A=\operatorname{diag}\left(\lambda_{1}, \ldots, \lambda_{k}\right)$ is diagonal. Then modulo suitable normalizing constants $\left(\frac{i}{2 \pi}\right)^{j}$, it is immediate that $c_{j}(A)$ is the $j^{\text {th }}$ elementary symmetric function of the $\lambda_{i}$ 's since

$$
\operatorname{det}(I+A)=\prod\left(1+\lambda_{j}\right)=1+s_{1}(\lambda)+\cdots+s_{k}(\lambda)
$$

If $P(\cdot)$ is any invariant polynomial, then $P(A)$ is a symmetric function of the $\lambda_{i}$ 's. The elementary symmetric functions form an algebra basis for the symmetric polynomials so there is a unique polynomial $Q$ so that

$$
P(A)=Q\left(c_{1}, \ldots, c_{k}\right)(A)
$$

i.e., we can decompose $P$ as a polynomial in the $c_{i}$ 's for diagonal matrices. Since $P$ is invariant, this is true for diagonalizable matrices. Since the diagonalizable matrices are dense and $P$ is continuous, this is true for all $A$. This proves:

Lemma 2.1.3. Let $P(A)$ be a invariant polynomial. There exists a unique polynomial $Q$ so that $P=Q\left(c_{1}(A), \ldots, c_{k}(A)\right)$.

It is clear that any polynomial in the $c_{k}$ 's is invariant, so this completely characterizes the ring of invariant polynomials; it is a free algebra in $k$ variables $\left\{c_{1}, \ldots, c_{k}\right\}$.

If $P(A)$ is homogeneous of degree $k$ and is defined on $k \times k$ matrices, we shall see that if $P(A) \neq 0$ as a polynomial, then there exists a holomorphic manifold $M$ so that if $T_{c}(M)$ is the holomorphic tangent space, then $\int_{M} P\left(T_{c}(M)\right) \neq 0$ where $\operatorname{dim} M=2 k$. This fact can be used to show that in general if $P(A) \neq 0$ as a polynomial, then there exists a manifold $M$ and a bundle $V$ over $M$ so $P(V) \neq 0$ in cohomology.

We can apply functorial constructions on connections. Define:

$$
\begin{array}{cll}
\nabla_{1} \oplus \nabla_{2} & \text { on } & C^{\infty}\left(V_{1} \oplus V_{2}\right) \\
\nabla_{1} \otimes \nabla_{2} & \text { on } & C^{\infty}\left(V_{1} \otimes V_{2}\right) \\
\nabla_{1}^{*} & \text { on } & C^{\infty}\left(V_{1}^{*}\right)
\end{array}
$$

by:

$$
\begin{aligned}
& \left(\nabla_{1} \oplus \nabla_{2}\right)\left(s_{1} \oplus s_{2}\right)=\left(\nabla_{1} s_{1}\right) \oplus\left(\nabla_{2} s_{2}\right) \\
& \quad \text { with } \omega=\omega_{1} \oplus \omega_{2} \text { and } \Omega=\Omega_{1} \oplus \Omega_{2} \\
& \left(\nabla_{1} \otimes \nabla_{2}\right)\left(s_{1} \otimes s_{2}\right)=\left(\nabla_{1} s_{1}\right) \otimes\left(\nabla_{2} s_{2}\right) \\
& \quad \text { with } \omega=\omega_{1} \otimes 1+1 \otimes \omega_{2} \text { and } \Omega=\Omega_{1} \otimes 1+1 \otimes \Omega_{2} \\
& \left(\nabla_{1} s_{1}, s_{1}^{*}\right)+\left(s_{1}, \nabla_{1}^{*} s_{1}^{*}\right)=d\left(s_{1}, s_{1}^{*}\right) \\
& \quad \text { relative to the dual frame } \omega=-\omega_{1}^{t} \text { and } \Omega=-\Omega_{1}^{t} .
\end{aligned}
$$

In a similar fashion we can define an induced connection on $\Lambda^{p}(V)$ (the bundle of $p$-forms) and $S^{p}(V)$ (the bundle of symmetric $p$-tensors). If $V$ has a given Hermitian fiber metric, the connection $\nabla$ is said to be unitary or Riemannian if $\left(\nabla s_{1}, s_{2}\right)+\left(s_{1}, \nabla s_{2}\right)=d\left(s_{1}, s_{2}\right)$. If we identify $V$ with $V^{*}$ using the metric, this simply means $\nabla=\bar{\nabla}^{*}$. This is equivalent to the condition that $\omega$ is a skew adjoint matrix of 1 -forms relative to a local orthonormal frame. We can always construct Riemannian connections locally by taking $\omega=0$ relative to a local orthonormal frame and then using a partition of unity to construct a global Riemannian connection.

If we restrict to Riemannian connections, then it is natural to consider polynomials $P(A)$ which are defined for skew-Hermitian matrices $A+A^{*}=$ 0 and which are invariant under the action of the unitary group. Exactly the same argument as that given in the proof of Lemma 2.1.3 shows that we can express such a $P$ in the form $P(A)=Q\left(c_{1}(A), \ldots, c_{k}(A)\right)$ so that the GL $(\cdot, \mathbf{C})$ and the $\mathrm{U}(\cdot)$ characteristic classes coincide. This is not true in the real category as we shall see; the Euler form is a characteristic form of the special orthogonal group which can not be defined as a characteristic form using the general linear group GL $(\cdot, \mathbf{R})$.

The Chern form and the Chern character satisfy certain identities with respect to functorial constructions.

Lemma 2.1.4.
(a)

$$
\begin{aligned}
c\left(V_{1} \oplus V_{2}\right) & =c\left(V_{1}\right) c\left(V_{2}\right) \\
c\left(V^{*}\right) & =1-c_{1}(V)+c_{2}(V)-\cdots+(-1)^{k} c_{k}(V)
\end{aligned}
$$

(b)

$$
\begin{aligned}
\operatorname{ch}\left(V_{1} \oplus V_{2}\right) & =\operatorname{ch}\left(V_{1}\right)+\operatorname{ch}\left(V_{2}\right) \\
\left.\operatorname{ch}\left(V_{1}\right) \oplus V_{2}\right) & =\operatorname{ch}\left(V_{1}\right) \operatorname{ch}\left(V_{2}\right)
\end{aligned}
$$

Proof: All the identities except the one involving $c\left(V^{*}\right)$ are immediate from the definition if we use the direct sum and product connections. If we fix a Hermitian structure on $V$, the identification of $V$ with $V^{*}$ is conjugate linear. The curvature of $\nabla^{*}$ is $-\Omega^{t}$ so we compute:

$$
c\left(V^{*}\right)=\operatorname{det}\left(I-\frac{i}{2 \pi} \Omega^{t}\right)=\operatorname{det}\left(I-\frac{i}{2 \pi} \Omega\right)
$$

which gives the desired identity.
If we choose a Riemannian connection on $V$, then $\Omega+\bar{\Omega}^{t}=0$. This immediately yields the identities:

$$
\overline{\operatorname{ch(V)}}=\operatorname{ch}(V) \quad \text { and } \quad \overline{c(V)}=c(V)
$$

so these are real cohomology classes. In fact the normalizing constants were chosen so $c_{k}(V)$ is an integral cohomology class-i.e., if $N_{2 k}$ is any oriented submanifold of dimension $2 k$, then $\int_{N_{2} k} c_{k}(V) \in \mathbf{Z}$. The $c h_{k}(V)$ are not integral for $k>1$, but they are rational cohomology classes. As we shall not need this fact, we omit the proof.

The characteristic classes give cohomological invariants of a vector bundle. We illustrate this by constructing certain examples over even dimensional spheres; these examples will play an important role in our later development of the Atiyah-Singer index theorem.

Definition. A set of matrices $\left\{e_{0}, \ldots, e_{m}\right\}$ are Clifford matrices if the $e_{i}$ are self-adjoint and satisfy the commutation relations $e_{i} e_{j}+e_{j} e_{i}=2 \delta_{i j}$ where $\delta_{i j}$ is the Kronecker symbol.

For example, if $m=2$ we can take:

$$
e_{0}=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right), \quad e_{1}=\left(\begin{array}{rr}
0 & 1 \\
1 & 0
\end{array}\right), \quad e_{2}=\left(\begin{array}{rr}
0 & i \\
-i & 0
\end{array}\right)
$$

to be the Dirac matrices. More generally, we can take the $e_{i}$ 's to be given by the spin representations. If $x \in \mathbf{R}^{m+1}$, we define:

$$
e(x)=\sum_{i} x_{i} e_{i} \quad \text { so } \quad e(x)^{2}=|x|^{2} I
$$

Conversely let $e(x)$ be a linear map from $\mathbf{R}^{m}$ to the set of self-adjoint matrices with $e(x)=|x|^{2} I$. If $\left\{v_{0}, \ldots, v_{m}\right\}$ is any orthonormal basis for $\mathbf{R}^{m},\left\{e\left(v_{0}\right), \ldots, e\left(v_{m}\right)\right\}$ forms a set of Clifford matrices.

If $x \in S^{m}$, we let $\Pi_{ \pm}(x)$ be the range of $\frac{1}{2}(1+e(x))=\pi_{ \pm}(x)$. This is the span of the $\pm 1$ eigenvectors of $e(x)$. If $e(x)$ is a $2 k \times 2 k$ matrix, then $\operatorname{dim} \Pi_{ \pm}(x)=k$. We have a decomposition $S^{m} \times \mathbf{C}^{2 k}=\Pi_{+} \oplus \Pi_{-}$. We project the flat connection on $S^{m} \times \mathbf{C}^{2 k}$ to the two sub-bundles to define connections $\nabla_{ \pm}$on $\Pi_{ \pm}$. If $e_{ \pm}^{0}$ is a local frame for $\Pi_{ \pm}\left(x_{0}\right)$, we define $e_{ \pm}(x)=\pi_{ \pm} e_{ \pm}^{0}$ as a frame in a neighborhood of $x_{0}$. We compute

$$
\nabla_{ \pm} e_{ \pm}=\pi_{ \pm} d \pi_{ \pm} e_{ \pm}^{0}, \quad \Omega_{ \pm} e_{ \pm}=\pi_{ \pm} d \pi_{ \pm} d \pi_{ \pm} e_{ \pm}^{0}
$$

Since $e_{ \pm}^{0}=e_{ \pm}\left(x_{0}\right)$, this yields the identity:

$$
\Omega_{ \pm}\left(x_{0}\right)=\pi_{ \pm} d \pi_{ \pm} d \pi_{ \pm}\left(x_{0}\right)
$$

Since $\Omega$ is tensorial, this holds for all $x$.

Let $m=2 j$ be even. We wish to compute $c h_{j}$. Suppose first $x_{0}=$ $(1,0, \ldots, 0)$ is the north pole of the sphere. Then:

$$
\begin{aligned}
\pi_{+}\left(x_{0}\right) & =\frac{1}{2}\left(1+e_{0}\right) \\
d \pi_{+}\left(x_{0}\right) & =\frac{1}{2} \sum_{i \geq 1} d x_{i} e_{i} \\
\Omega_{+}\left(x_{0}\right) & =\frac{1}{2}\left(1+e_{0}\right)\left(\frac{1}{2} \sum_{i \geq 1} d x_{i} e_{i}\right)^{2} \\
\Omega_{+}\left(x_{0}\right)^{j} & =\frac{1}{2}\left(1+e_{0}\right)\left(\frac{1}{2} \sum_{i \geq 1} d x_{i} e_{i}\right)^{2 j} \\
& =2^{-m-1} m!\left(1+e_{0}\right)\left(e_{1} \ldots e_{m}\right)\left(d x_{1} \wedge \cdots \wedge d x_{m}\right) .
\end{aligned}
$$

The volume form at $x_{0}$ is $d x_{1} \wedge \cdots \wedge d x_{m}$. Since $e_{1}$ anti-commutes with the matrix $e_{1} \ldots e_{m}$, this matrix has trace 0 so we compute:

$$
c h_{j}\left(\Omega_{+}\right)\left(x_{0}\right)=\left(\frac{i}{2 \pi}\right)^{j} 2^{-m-1} m!\operatorname{Tr}\left(e_{0} \ldots e_{m}\right) \operatorname{dvol}\left(x_{0}\right) / j!
$$

A similar computation shows this is true at any point $x_{0}$ of $S^{m}$ so that:

$$
\int_{S^{m}} c h_{j}\left(\Pi_{+}\right)=\left(\frac{i}{2 \pi}\right)^{j} 2^{-m-1} m!\operatorname{Tr}\left(e_{0} \ldots e_{m}\right) \operatorname{vol}\left(s^{m}\right) / j!
$$

Since the volume of $S^{m}$ is $j!2^{m+1} \pi^{j} / m$ ! we conclude:
Lemma 2.1.5. Let $e(x)$ be a linear map from $\mathbf{R}^{m+1}$ to the set of selfadjoint matrices. We suppose $e(x)^{2}=|x|^{2} I$ and define bundles $\Pi_{ \pm}(x)$ over $S^{m}$ corresponding to the $\pm 1$ eigenvalues of $e$. Let $m=2 j$ be even, then:

$$
\int_{S^{m}} c h_{j}\left(\Pi_{+}\right)=i^{j} 2^{-j} \operatorname{Tr}\left(e_{0} \ldots e_{m}\right)
$$

In particular, this is always an integer as we shall see later when we consider the spin complex. If

$$
e_{0}=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right), \quad e_{1}=\left(\begin{array}{rr}
0 & 1 \\
1 & 0
\end{array}\right), \quad e_{2}=\left(\begin{array}{rr}
0 & i \\
-i & 0
\end{array}\right)
$$

then $\operatorname{Tr}\left(e_{0} e_{1} e_{2}\right)=-2 i$ so $\int_{S^{2}} c h_{1}\left(\Pi_{+}\right)=1$ which shows in particular $\Pi_{+}$ is a non-trivial line bundle over $S^{2}$.

There are other characteristic classes arising in the index theorem. These are most conveniently discussed using generating functions. Let

$$
x_{j}=\frac{i}{2 \pi} \lambda_{j}
$$

be the normalized eigenvalues of the matrix $A$. We define:

$$
c(A)=\prod_{j=1}^{k}\left(1+x_{j}\right) \quad \text { and } \quad \operatorname{ch}(A)=\sum_{j=1}^{k} e^{x_{j}} .
$$

If $P(x)$ is a symmetric polynomial, we express $P(x)=Q\left(c_{1}(x), \ldots, c_{k}(x)\right)$ to show $P(A)$ is a polynomial in the components of $A$. More generally if $P$ is analytic, we first express $P$ in a formal power series and then collect homogeneous terms to define $P(A)$. We define:

Todd class: $\operatorname{Td}(A)=\prod_{j=1}^{k} \frac{x_{j}}{1-e^{-x_{j}}}$

$$
=1+\frac{1}{2} c_{1}(A)+\frac{1}{12}\left(c_{1}^{2}+c_{2}\right)(A)+\frac{1}{24} c_{1}(A) c_{2}(A)+\cdots
$$

The Todd class appears in the Riemann-Roch theorem. It is clear that it is multiplicative with respect to direct sum- $T d\left(V_{1} \oplus V_{2}\right)=T d\left(V_{1}\right) T d\left(V_{2}\right)$. The Hirzebruch L-polynomial and the $\hat{A}$ polynomial will be discussed in the next subsection. These are real characteristic classes which also are defined by generating functions.

If $V$ is a bundle of dimension $k$, then $V \oplus 1$ will be a bundle of dimension $k+1$. It is clear $c(V \oplus 1)=c(V)$ and $T d(V \oplus 1)=T d(V)$ so these are stable characteristic classes. $\operatorname{ch}(V)$ on the other hand is not a stable characteristic class since $c h_{0}(V)=\operatorname{dim} V$ depends explicitly on the dimension of $V$ and changes if we alter the dimension of $V$ by adding a trivial bundle.

### 2.2. Characteristic Classes of a Real Vector Bundle. Pontrjagin and Euler Classes.

Let $V$ be a real vector bundle of dimension $k$ and let $V_{c}=V \otimes \mathbf{C}$ denote the complexification of $V$. We place a real fiber metric on $V$ to reduce the structure group from $\operatorname{GL}(k, \mathbf{R})$ to the orthogonal group $\mathrm{O}(k)$. We restrict henceforth to Riemannian connections on $V$, and to local orthonormal frames. Under these assumptions, the curvature is a skew-symmetric matrix of 2-forms:

$$
\Omega+\Omega^{t}=0
$$

Since $V_{c}$ arises from a real vector bundle, the natural isomorphism of $V$ with $V^{*}$ defined by the metric induces a complex linear isomorphism of $V_{c}$ with $V_{c}^{*}$ so $c_{j}\left(V_{c}\right)=0$ for $j$ odd by Lemma 2.1.4. Expressed locally,

$$
\operatorname{det}\left(I+\frac{i}{2 \pi} A\right)=\operatorname{det}\left(I+\frac{i}{2 \pi} A^{t}\right)=\operatorname{det}\left(I-\frac{i}{2 \pi} A\right)
$$

if $A+A^{t}=0$ so $c(A)$ is an even polynomial in $A$. To avoid factors of $i$ we define the Pontrjagin form:

$$
p(A)=\operatorname{det}\left(I+\frac{1}{2 \pi} A\right)=1+p_{1}(A)+p_{2}(A)+\cdots
$$

where $p_{j}(A)$ is homogeneous of order $2 j$ in the components of $A$; the corresponding characteristic class $p_{j}(V) \in H^{4 j}(M ; \mathbf{R})$. It is immediate that:

$$
p_{j}(V)=(-1)^{j} c_{2 j}\left(V_{c}\right)
$$

where the factor of $(-1)^{j}$ comes form the missing factors of $i$.
The set of skew-symmetric matrices is the Lie algebra of $\mathrm{O}(k)$. Let $P(A)$ be an invariant polynomial under the action of $\mathrm{O}(k)$. We define $P(\Omega)$ for $\nabla$ a Riemannian connection exactly as in the previous subsections. Then the analogue of Lemma 2.1.3 is:

Lemma 2.2.1. Let $P(A)$ be a polynomial in the components of a matrix $A$. Suppose $P(A)=P\left(g A g^{-1}\right)$ for every skew-symmetric $A$ and for every $g \in \mathrm{O}(k)$. Then there exists a polynomial $Q\left(p_{1}, \ldots\right)$ so $P(A)=$ $Q\left(p_{1}(A), \ldots\right)$ for every skew-symmetric $A$.

Proof: It is important to note that we are not asserting that we have $P(A)=Q\left(p_{1}(A), \ldots\right)$ for every matrix $A$, but only for skew-symmetric $A$. For example, $P(A)=\operatorname{Tr}(A)$ vanishes for skew symmetric $A$ but does not vanish in general.

It is not possible in general to diagonalize a skew-symmetric real matrix. We can, however, put it in block diagonal form:

$$
A=\left(\begin{array}{ccccc}
0 & -\lambda_{1} & 0 & 0 & \ldots \\
\lambda_{1} & 0 & 0 & 0 & \ldots \\
0 & 0 & 0 & -\lambda_{2} & \ldots \\
0 & 0 & \lambda_{2} & 0 & \ldots
\end{array}\right) .
$$

If $k$ is odd, then the last block will be a $1 \times 1$ block with a zero in it. We let $x_{j}=-\lambda_{j} / 2 \pi$; the sign convention is chosen to make the Euler form have the right sign. Then:

$$
p(A)=\prod_{j}\left(1+x_{j}^{2}\right)
$$

where the product ranges from 1 through $\left[\frac{k}{2}\right]$.
If $P(A)$ is any invariant polynomial, then $P$ is a symmetric function in the $\left\{x_{j}\right\}$. By conjugating $A$ by an element of $\mathrm{O}(k)$, we can replace any $x_{j}$ by $-x_{j}$ so $P$ is a symmetric function of the $\left\{x_{j}^{2}\right\}$. The remainder of the proof of Lemma 2.2.1 is the same; we simply express $P(A)$ in terms of the elementary symmetric functions of the $\left\{x_{j}^{2}\right\}$.

Just as in the complex case, it is convenient to define additional characteristic classes using generating functions. The functions:

$$
z / \tanh z \quad \text { and } \quad z\{2 \sinh (z / 2)\}
$$

are both even functions of the parameter $z$. We define:
The Hirzebruch L-polynomial:

$$
\begin{aligned}
L(x) & =\prod_{j} \frac{x_{j}}{\tanh x_{j}} \\
& =1+\frac{1}{3} p_{1}+\frac{1}{45}\left(7 p_{2}-p_{1}^{2}\right) \cdots
\end{aligned}
$$

The $\hat{A}$ (A-roof) genus:

$$
\begin{aligned}
A(x) & =\prod_{j} \frac{x_{j}}{2 \sinh \left(x_{j} / 2\right)} \\
& =1-\frac{1}{24} p_{1}+\frac{1}{5760}\left(7 p_{1}^{2}-4 p_{2}\right)+\cdots
\end{aligned}
$$

These characteristic classes appear in the formula for the signature and spin complexes.

For both the real and complex case, the characteristic ring is a pure polynomial algebra without relations. Increasing the dimension $k$ just adds generators to the ring. In the complex case, the generators are the Chern classes $\left\{1, c_{1}, \ldots, c_{k}\right\}$. In the real case, the generators are the Pontrjagin classes $\left\{1, p_{1}, \ldots, p_{[k / 2]}\right\}$ where $[k / 2]$ is the greatest integer in $k / 2$. There is one final structure group which will be of interest to us.

Let $V$ be a vector bundle of real dimension $k . V$ is orientable if we can choose a fiber orientation consistently. This is equivalent to assuming that the real line bundle $\Lambda^{k}(V)$ is trivial. We choose a fiber metric for $V$ and an orientation for $V$ and restrict attention to oriented orthonormal frames. This restricts the structure group from $\mathrm{O}(k)$ to the special orthogonal group $\mathrm{SO}(k)$.

If $k$ is odd, no new characteristic classes result from the restriction of the fiber group. We use the final $1 \times 1$ block of 0 in the representation of $A$ to replace any $\lambda_{i}$ by $-\lambda_{i}$ by conjugation by an element of $\mathrm{SO}(k)$. If $n$ is even, however, we cannot do this as we would be conjugating by a orientation reversing matrix.

Let $P(A)$ be a polymonial in the components of $A$. Suppose $P(A)=$ $P\left(g A g^{-1}\right)$ for every skew-symmetric $A$ and every $g \in \operatorname{SO}(k)$. Let $k=2 \bar{k}$ be even and let $P(A)=P\left(x_{1}, \ldots, x_{\bar{k}}\right)$. Fix $g_{0} \in \mathrm{O}(k)-\mathrm{SO}(k)$ and define:

$$
P_{0}(A)=\frac{1}{2}\left(P(A)+P\left(g_{0} A g_{0}^{-1}\right)\right) \quad \text { and } \quad P_{1}(A)=\frac{1}{2}\left(P(A)-P\left(g_{0} A g_{0}^{-1}\right)\right)
$$

Both $P_{0}$ and $P_{1}$ are $\mathrm{SO}(k)$ invariant. $P_{0}$ is $\mathrm{O}(k)$ invariant while $P_{1}$ changes sign under the action of an element of $\mathrm{O}(k)-\mathrm{SO}(k)$.

We can replace $x_{1}$ by $-x_{1}$ by conjugating by a suitable orientation reversing map. This shows:

$$
P_{1}\left(x_{1}, x_{2}, \ldots\right)=-P\left(-x_{1}, x_{2}, \ldots\right)
$$

so that $x_{1}$ must divide every monomial of $P_{0}$. By symmetry, $x_{j}$ divides every monomial of $P_{1}$ for $1 \leq j \leq \bar{k}$ so we can express:

$$
P_{1}(A)=x_{1} \ldots x_{\bar{k}} P_{1}^{\prime}(A)
$$

where $P_{1}^{\prime}(A)$ is now invariant under the action of $\mathrm{O}(k)$. Since $P_{1}^{\prime}$ is polynomial in the $\left\{x_{j}^{2}\right\}$, we conclude that both $P_{0}$ and $P_{1}^{\prime}$ can be represented as polynomials in the Pontrjagin classes. We define:

$$
e(A)=x_{1} \ldots x_{k} \quad \text { so } \quad e(A)^{2}=\operatorname{det}(A)=p_{\bar{k}}(A)=\prod_{j} x_{j}^{2}
$$

and decompose:

$$
P=P_{0}+e(A) P_{1}^{\prime}
$$

$e(A)$ is a square root of the determinant of $A$.
It is not, of course, immediate that $e(A)$ is a polynomial in the components of $A$. Let $\left\{v_{i}\right\}$ be an oriented orthonormal basis for $\mathbf{R}^{k}$. We let $A v_{i}=\sum_{j} A_{i j} v_{j}$ and define

$$
\omega(A)=\frac{1}{2 \pi} \sum_{i<j} A_{i j} v_{i} \wedge v_{j} \in \Lambda^{2}\left(\mathbf{R}^{k}\right)
$$

We let $v_{1} \wedge \cdots \wedge v_{k}$ be the orientation class of $\mathbf{R}^{k}$ and define:

$$
e(A)=\left(\omega(A)^{\bar{k}}, v_{1} \wedge \cdots \wedge v_{k}\right) / k!
$$

where (, ) denotes the natural inner product on $\Lambda^{k} \mathbf{R}^{k} \simeq \mathbf{R}$. It is clear from the definition that $e(A)$ is invariant under the action of $\mathrm{SO}(k)$ since $\omega(A)$ is invariantly defined for skew-symmetric matrices $A$. It is clear that $e(A)$ is polynomial in the components of $A$. If we choose a block basis so that:

$$
A v_{1}=\lambda_{1} v_{2}, \quad A v_{2}=-\lambda_{1} v_{1}, \quad A v_{3}=\lambda_{2} v_{4}, \quad A v_{4}=-\lambda_{2} v_{3}, \ldots
$$

then we compute:

$$
\begin{aligned}
\omega(A) & =\left\{-\lambda_{1} v_{1} \wedge v_{2}-\lambda_{2} v_{3} \wedge v_{4} \ldots\right\} / 2 \pi=x_{1} v_{1} \wedge v_{2}+x_{2} v_{3} \wedge v_{4} \cdots \\
e(A) & =x_{1} x_{2} \ldots
\end{aligned}
$$

This new characteristic class is called the Euler class. While the Pontrjagin classes can be computed from the curvature of an arbitrary connection, the Euler class can only be computed from the curvature of a Riemannian connection. If $\Omega_{i j}$ are the matrix components of the curvature of $V$ relative to some oriented orthomormal basis, then:

$$
e(\Omega)=(-4 \pi)^{-\bar{k}} / k!\sum_{\rho} \operatorname{sign}(\rho) \Omega_{\rho(1) \rho(2)} \ldots \Omega_{\rho(k-1) \rho(k)} \in \Lambda^{k}(M)
$$

for $2 \bar{k}=k=\operatorname{dim}(V)$. The sum ranges over all possible permutations $\rho$.
We define $e(V)=0$ if $\operatorname{dim}(V)$ is odd. It is immediate that $e\left(V_{1} \oplus V_{2}\right)=$ $e\left(V_{1}\right) e\left(V_{2}\right)$ if we give the natural orientation and fiber metric to the direct sum and if we use the direct sum connection.

We illustrate this formula (and check that we have the correct normalizing constants) by studying the following simple example. Let $m=2$ and let $M=S^{2}$ be the unit sphere. We calculate $e\left(T S^{2}\right)$. Parametrize $M$ using spherical coordinates $f(u, v)=(\cos u \sin v, \sin u \sin v, \cos v)$ for $0 \leq u \leq 2 \pi$ and $0 \leq v \leq \pi$. Define a local orthonormal frame for $T\left(\mathbf{R}^{3}\right)$ over $S^{2}$ by:

$$
\begin{aligned}
e_{1} & =(\sin v)^{-1} \partial / \partial u
\end{aligned}=(-\sin u, \cos u, 0), ~ \begin{aligned}
\partial / \partial v & =(\cos u \cos v, \sin u \cos v,-\sin v) \\
e_{2} & =\quad N \\
e_{3} & =(\cos u \sin v, \sin u \sin v, \cos v) .
\end{aligned}
$$

The Euclidean connection $\widetilde{\nabla}$ is easily computed to be:

$$
\begin{aligned}
& \widetilde{\nabla}_{e_{1}} e_{1}=(\sin v)^{-1}(-\cos u,-\sin u, 0)=-(\cos v / \sin v) e_{2}-e_{3} \\
& \widetilde{\nabla}_{e_{1}} e_{2}=(\cos v / \sin v)(-\sin u, \cos u, 0)=(\cos v / \sin v) e_{1} \\
& \widetilde{\nabla}_{e_{2}} e_{1}=0 \\
& \widetilde{\nabla}_{e_{2}} e_{2}=(-\cos u \sin v,-\sin u \sin v,-\cos v)=-e_{3} .
\end{aligned}
$$

Covariant differentiation $\nabla$ on the sphere is given by projecting back to $T\left(S^{2}\right)$ so that:

$$
\begin{array}{ll}
\nabla_{e_{1}} e_{1}=(-\cos v / \sin v) e_{2} & \nabla_{e_{1} e_{2}}=(\cos v / \sin v) e_{1} \\
\nabla_{e_{2}} e_{1}=0 & \nabla_{e_{2}} e_{2}=0
\end{array}
$$

and the connection 1-form is given by:

$$
\begin{array}{ll}
\nabla e_{1}=(-\cos v / \sin v) e^{1} \otimes e_{2} & \text { so } \omega_{11}=0 \text { and } \omega_{12}=-(\cos v / \sin v) e^{1} \\
\nabla e_{2}=(\cos v / \sin v) e^{1} \otimes e_{1} & \text { so } \omega_{22}=0 \text { and } \omega_{21}=(\cos v / \sin v) e^{1}
\end{array}
$$

As $d u=e^{1} / \sin v$ we compute $\omega_{12}=-\cos v d u$ so $\Omega_{12}=\sin v d v \wedge d u=$ $-e^{1} \wedge e^{2}$ and $\Omega_{21}=e^{1} \wedge e^{2}$. From this we calculate that:

$$
e(\Omega)=-\frac{1}{4 \pi}\left(\Omega_{12}-\Omega_{21}\right)=\frac{e^{1} \wedge e^{2}}{2 \pi}
$$

and consequently $\int_{S^{2}} E_{2}=\operatorname{vol}\left(S^{2}\right) / 2 \pi=2=\chi\left(S^{2}\right)$.
There is a natural relation between the Euler form and $c_{\bar{k}}$. Let $V$ be a complex vector space of dimension $\bar{k}$ and let $V_{r}$ be the underlying real vector space of dimention $2 \bar{k}=k$. If $V$ has a Hermitian inner product, then $V_{r}$ inherits a natural real inner product. $V_{r}$ also inherits a natural orientation from the complex structure on $V$. If $\left\{e_{j}\right\}$ is a unitary basis for $V$, then $\left\{e_{1}, i e_{1}, e_{2}, i e_{2}, \ldots\right\}$ is an oriented orthonormal basis for $V_{r}$. Let $A$ be a skew-Hermitian matrix on $V$. The restriction of $A$ to $V_{r}$ defines a skewsymmetric matrix $A_{r}$ on $V_{r}$. We choose a basis $\left\{e_{j}\right\}$ for $V$ so $A e_{j}=i \lambda_{j} e_{j}$. Then:

$$
x_{j}=-\lambda_{j} / 2 \pi \quad \text { and } \quad c_{\bar{k}}(A)=x_{1} \ldots x_{\bar{k}}
$$

defines the $\bar{k}^{\text {th }}$ Chern class. If $v_{1}=e_{1}, v_{2}=i e_{1}, v_{3}=e_{2}, v_{4}=i e_{2}, \ldots$ then:

$$
A_{r} v_{1}=\lambda_{1} v_{2}, \quad A_{r} v_{2}=-\lambda_{1} v_{1}, \quad A_{r} v_{3}=\lambda_{2} v_{4}, \quad A_{r} v_{4}=-\lambda_{2} v_{3}, \ldots
$$

so that

$$
e\left(A_{r}\right)=c_{\bar{k}}(A)
$$

We summarize these results as follows:

Lemma 2.2.2. Let $P(A)$ be a polynomial with $P(A)=P\left(g A g^{-1}\right)$ for every skew symmetric $A$ and every $g \in \operatorname{SO}(k)$.
(a) If $k$ is odd, then $P(A)$ is invariant under $\mathrm{O}(k)$ and is expressible in terms of Pontrjagin classes.
(b) If $k=2 \bar{k}$ is even, then we can decompose $P(A)=P_{0}(A)+e(A) P_{1}(A)$ where $P_{i}(A)$ are $\mathrm{O}(k)$ invariant and are expressible in terms of Pontrjagin classes.
(c) $e(A)$ is defined by:

$$
e(A)=(-4 \pi)^{-\bar{k}} \sum_{\rho} \operatorname{sign}(\rho) A_{\rho(1) \rho(2)} \ldots A_{\rho(k-1) \rho(k)} / k!
$$

This satisfies the identity $e(A)^{2}=p_{\bar{k}}(A)$. (We define $e(A)=0$ if $k$ is odd).
This completely describes the characteristic ring. We emphasize the conclusions are only applicable to skew-symmetric real matrices $A$. We have proved that $e(A)$ has the following functorial properties:

Lemma 2.2.3.
(a) If we take the direct sum connection and metric, then $e\left(V_{1} \oplus V_{2}\right)=$ $e\left(V_{1}\right) e\left(V_{2}\right)$.
(b) If we take the metric and connection of $V_{r}$ obtained by forgetting the complex structure on a complex bundle $V$ then $e\left(V_{r}\right)=c_{\bar{k}}(V)$.

This lemma establishes that the top dimensional Chern class $c_{\bar{k}}$ does not really depend on having a complex structure but only depends on having an orientation on the underlying real vector bundle. The choice of the sign in computing $\operatorname{det}(A)^{1 / 2}$ is, of course, motivated by this normalization.

### 2.3. Characteristic Classes of Complex Projective Space.

So far, we have only discussed covariant differentiation from the point of view of the total covariant derivative. At this stage, it is convenient to introduce covariant differentiation along a direction. Let $X \in T(M)$ be a tangent vector and let $s \in C^{\infty}(V)$ be a smooth section. We define $\nabla_{X} s \in C^{\infty}(V)$ by:

$$
\nabla_{X} s=X \cdot \nabla s
$$

where "." denotes the natural pairing from $T(M) \otimes T^{*}(M) \otimes V \rightarrow V$. Let $[X, Y]=X Y-Y X$ denote the Lie-bracket of vector fields. We define:

$$
\Omega(X, Y)=\nabla_{X} \nabla_{Y}-\nabla_{Y} \nabla_{X}-\nabla_{[X, Y]}
$$

and compute that:

$$
\begin{aligned}
& \Omega(f X, Y) s=\Omega(X, f Y) s=\Omega(X, Y) f s=f \Omega(X, Y) s \\
& \quad \text { for } \quad f \in C^{\infty}(M), s \in C^{\infty}(V), \text { and } X, Y \in C^{\infty}(T M) .
\end{aligned}
$$

If $\left\{e_{i}\right\}$ is a local frame for $T(M)$, let $\left\{e^{i}\right\}$ be the dual frame for $T^{*}(M)$. It is immediate that:

$$
\nabla s=\sum_{i} e^{i} \otimes \nabla_{e_{i}}(s)
$$

so we can recover the total covariant derivative from the directional derivatives. We can also recover the curvature:

$$
\Omega s=\sum_{i<j} e^{i} \wedge e^{j} \otimes \Omega\left(e_{i}, e_{j}\right) s
$$

If $V=T(M)$, there is a special connection called the Levi-Cività connection on $T(M)$. It is the unique Riemannian connection which is torsion free-i.e.,

$$
\begin{aligned}
\left(\nabla_{X} Y, Z\right)+\left(Y, \nabla_{X} Z\right) & =X(Y, Z) \\
\nabla_{X} Y-\nabla_{Y} X-[X, Y] & =0
\end{aligned}
$$

Let $x=\left(x_{1}, \ldots, x_{m}\right)$ be a system of local coordinates on $M$ and let $\left\{\partial / \partial x_{i}\right\}$ be the coordinate frame for $T(M)$. The metric tensor is given by:

$$
d s^{2}=\sum_{i, j} g_{i j} d x^{i} \otimes d x^{j} \quad \text { for } \quad g_{i j}=\left(\partial / \partial x_{i}, \partial / \partial x_{j}\right)
$$

We introduce the Christoffel symbols $\Gamma_{i j}^{k}$ and $\Gamma_{i j k}$ by:

$$
\begin{aligned}
& \nabla_{\partial / \partial x_{i}} \partial / \partial x_{j}=\sum_{k} \Gamma_{i j}^{k} \partial / \partial x_{k} \\
& \left(\nabla_{\partial / \partial x_{i}} \partial / \partial x_{j}, \partial / \partial x_{k}\right)=\Gamma_{i j k}
\end{aligned}
$$

They are related by the formula:

$$
\Gamma_{i j k}=\sum_{l} \Gamma_{i j}^{l} g_{l k}, \quad \Gamma_{i j}^{k}=\sum_{l} \Gamma_{i j l} g^{l k}
$$

where $g^{l k}=\left(d x_{l}, d x_{k}\right)$ is the inverse of the matrix $g_{i j}$. It is not difficult to compute that:

$$
\Gamma_{i j k}=\frac{1}{2}\left\{g_{j k / i}+g_{i k / j}-g_{i j / k}\right\}
$$

where we use the notation "/" to denote (multiple) partial differentiation.
The complete curvature tensor is defined by:

$$
\Omega\left(\partial / \partial x_{i}, \partial / \partial x_{j}\right) \partial / \partial x_{k}=\sum_{l} R_{i j k}^{l} \partial / \partial x_{l}
$$

or equivalently if we lower indices:

$$
\left(\Omega\left(\partial / \partial x_{i}, \partial / \partial x_{j}\right) \partial / \partial x_{k}, \partial / \partial x_{l}\right)=R_{i j k l} .
$$

The expression of the curvature tensor in terms of the derivatives of the metric is very complicated in general. By making a linear change of coordinates we can always normalize the metric so $g_{i j}\left(x_{0}\right)=\delta_{i j}$ is the Kronecker symbol. Similarly, by making a quadratic change of coordinates, we can further normalize the metric so $g_{i j / k}\left(x_{0}\right)=0$. Relative to such a choice of coordinates:

$$
\Gamma_{i j}^{k}\left(x_{0}\right)=\Gamma_{i j k}\left(x_{0}\right)=0
$$

and

$$
R_{i j k l}\left(x_{0}\right)=R_{i j k}^{l}\left(x_{0}\right)=\frac{1}{2}\left\{g_{j l / i k}+g_{i k / l j}-g_{j k / i l}-g_{i l / j k}\right\}\left(x_{0}\right)
$$

At any other point, of course, the curvature tensor is not as simply expressed. In general it is linear in the second derivatives of the metric, quadratic in the first derivatives of the metric with coefficients which depend smoothly on the $g_{i j}$ 's.

We choose a local orthonormal frame $\left\{e_{i}\right\}$ for $T^{*}(M)$. Let $m=2 \bar{m}$ be even and let orn $=* \cdot 1=$ dvol be the oriented volume. Let

$$
e=E_{m}(g) \text { orn }
$$

be the Euler form. If we change the choice of the local orientation, then $e$ and orn both change signs so $E_{m}(g)$ is scalar invariant of the metric. In terms of the curvature, if $2 \bar{m}=m$, then:

$$
\begin{aligned}
c(m) & =(-1)^{\bar{m}}(8 \pi)^{-\bar{m}} \frac{1}{(\bar{m})!} \\
E_{m} & =c(m) \sum_{\rho, \tau} \operatorname{sign}(\rho) \operatorname{sign}(\tau) R_{\rho(1) \rho(2) \tau(1) \tau(2)} \ldots R_{\rho(m-1) \rho(m) \tau(m-1) \tau(m)}
\end{aligned}
$$

where the sum ranges over all permutations $\rho, \tau$ of the integers 1 through $m$. For example:

$$
\begin{aligned}
& E_{2}=-(2 \pi)^{-1} R_{1212} \\
& E_{4}=(2 \pi)^{-2} \sum_{i, j, k, l}\left\{R_{i j i j} R_{k l k l}+R_{i j k l} R_{i j k l}-4 R_{i j i k} R_{l j l k}\right\} / 8
\end{aligned}
$$

Let dvol be the Riemannian measure on $M$. If $M$ is oriented, then

$$
\int_{M} E_{m} \mathrm{dvol}=\int_{M} e
$$

is independent of the orientation of $M$ and of the metric. If $M$ is not orientable, we pass to the double cover to see $\int_{M} E_{m}$ dvol is a topological invariant of the manifold $M$. We shall prove later this integral is the Euler characteristic $\chi(M)$ but for the moment simply note it is not dependent upon a choice of orientation of $M$ and is in fact defined even if $M$ is not orientable.

It is worth computing an example. We let $S^{2}$ be the unit sphere in $\mathbf{R}^{3}$. Since this is homogeneous, $E_{2}$ is constant on $S^{2}$. We compute $E_{2}$ at the north pole $(0,0,1)$ and parametrize $S^{2}$ by $\left(u, v,\left(1-u^{2}-v^{2}\right)^{1 / 2}\right)$. Then

$$
\begin{aligned}
\partial / \partial u & =\left(1,0,-u /\left(1-u^{2}-v^{2}\right)^{1 / 2}\right), \quad \partial / \partial v=\left(0,1,-v /\left(1-u^{2}-v^{2}\right)^{1 / 2}\right) \\
g_{11} & =1+u^{2} /\left(1-u^{2}-v^{2}\right) \\
g_{22} & =1+v^{2} /\left(1-u^{2}-v^{2}\right) \\
g_{12} & =u v /\left(1-u^{2}-v^{2}\right)
\end{aligned}
$$

It is clear $g_{i j}(0)=\delta_{i j}$ and $g_{i j / k}(0)=0$. Therefore at $u=v=0$,

$$
\begin{aligned}
E_{2} & =-(2 \pi)^{-1} R_{1212}=-(2 \pi)^{-1}\left\{g_{11 / 22}+g_{22 / 11}-2 g_{12 / 12}\right\} / 2 \\
& =-(2 \pi)^{-1}\{0+0-2\} / 2=(2 \pi)^{-1}
\end{aligned}
$$

so that

$$
\int_{S^{2}} E_{2} \mathrm{dvol}=(4 \pi)(2 \pi)^{-1}=2=\chi\left(S^{2}\right)
$$

More generally, we can let $M=S^{2} \times \cdots \times S^{2}$ where we have $\bar{m}$ factors and $2 \bar{m}=m$. Since $e\left(V_{1} \oplus V_{2}\right)=e\left(V_{1}\right) e\left(V_{2}\right)$, when we take the product metric we have $E_{m}=\left(E_{2}\right)^{\bar{m}}=(2 \pi)^{-\bar{m}}$ for this example. Therefore:

$$
\int_{S^{2} \times \cdots \times S^{2}} E_{m} \mathrm{dvol}=2^{\bar{m}}=\chi\left(S^{2} \times \cdots \times S^{2}\right)
$$

We will use this example later in this chapter to prove $\int_{M} E_{m}=\chi(M)$ in general. The natural examples for studying the Euler class are products
of two dimensional spheres. Unfortunately, the Pontrjagin classes vanish identically on products of spheres so we must find other examples.

It is convenient at this point to discuss the holomorphic category. A manifold $M$ of real dimension $m=2 \bar{m}$ is said to be holomorphic if we have local coordinate charts $z=\left(z_{1}, \ldots, z_{\bar{m}}\right): M \rightarrow \mathbf{C}^{\bar{m}}$ such that two charts are related by holomorphic changes of coordinates. We expand $z_{j}=x_{j}+i y_{j}$ and define:

$$
\begin{aligned}
\partial / \partial z_{j} & =\frac{1}{2}\left(\partial / \partial x_{j}-i \partial / \partial y_{j}\right), & \partial / \partial \bar{z}_{j} & =\frac{1}{2}\left(\partial / \partial x_{j}+i \partial / \partial y_{j}\right) \\
d z_{j} & =d x_{j}+i d y_{j}, & d \bar{z}_{j} & =d x_{j}+i d y_{j} .
\end{aligned}
$$

We complexify $T(M)$ and $T^{*}(M)$ and define:

$$
\begin{aligned}
T_{c}(M) & =\operatorname{span}\left\{\partial / \partial z_{j}\right\} \\
\Lambda^{1,0}(M) & =\operatorname{span}\left\{d z_{j}\right\} \\
\Lambda^{0,1}(M) & =\operatorname{span}\left\{d \bar{z}_{j}\right\}
\end{aligned}
$$

Then

$$
\begin{aligned}
\Lambda^{1}(M) & =\Lambda^{1,0}(M) \oplus \Lambda^{0,1}(M) \\
T_{c}(M)^{*} & =\Lambda^{1,0}(M)
\end{aligned}
$$

As complex bundles, $T_{c}(M) \simeq \Lambda^{0,1}(M)$. Define:

$$
\begin{aligned}
& \partial(f)=\sum_{j} \partial f / \partial z_{j} d z^{j}: C^{\infty}(M) \rightarrow C^{\infty}\left(\Lambda^{1,0}(M)\right) \\
& \bar{\partial}(f)=\sum_{j} \partial f / \partial \bar{z}_{j} d \bar{z}^{j}: C^{\infty}(M) \rightarrow C^{\infty}\left(\Lambda^{0,1}(M)\right)
\end{aligned}
$$

then the Cauchy-Riemann equations show $f$ is holomorphic if and only if $\bar{\partial} f=0$.

This decomposes $d=\partial+\bar{\partial}$ on functions. More generally, we define:

$$
\begin{aligned}
\Lambda^{p, q} & =\operatorname{span}\left\{d z_{i_{1}} \wedge \cdots \wedge d z_{i_{p}} \wedge d \bar{z}_{j_{1}} \wedge \cdots \wedge d \bar{z}_{j_{q}}\right\} \\
\Lambda^{n} & =\bigoplus_{p+q=n} \Lambda^{p, q}
\end{aligned}
$$

This spanning set of $\Lambda^{p, q}$ is closed. We decompose $d=\partial+\bar{\partial}$ where

$$
\partial: C^{\infty}\left(\Lambda^{p, q}\right) \rightarrow C^{\infty}\left(\Lambda^{p+1, q}\right) \quad \text { and } \quad \bar{\partial}: C^{\infty}\left(\Lambda^{p, q}\right) \rightarrow C^{\infty}\left(\Lambda^{p, q+1}\right)
$$

so that $\partial \partial=\bar{\partial} \bar{\partial}=\bar{\partial} \partial+\partial \bar{\partial}=0$. These operators and bundles are all invariantly defined independent of the particular holomorphic coordinate system chosen.

A bundle $V$ over $M$ is said to be holomorphic if we can cover $M$ by holomorphic charts $U_{\alpha}$ and find frames $s_{\alpha}$ over each $U_{\alpha}$ so that on $U_{\alpha} \cap U_{\beta}$ the transition functions $s_{\alpha}=f_{\alpha \beta} s_{\beta}$ define holomorphic maps to GL $(n, \mathbf{C})$. For such a bundle, we say a local section $s$ is holomorphic if $s=\sum_{\nu} a_{\alpha}^{\nu} s_{\alpha}^{\nu}$ for holomorphic functions $a_{\alpha}^{\nu}$. For example, $T_{c}(M)$ and $\Lambda^{p, 0}(M)$ are holomorphic bundles over $M ; \Lambda^{0,1}$ is anti-holomorphic.

If $V$ is holomorphic, we use a partition of unity to construct a Hermitian fiber metric $h$ on $V$. Let $h_{\alpha}=\left(s_{\alpha}, s_{\alpha}\right)$ define the metric locally; this is a positive definite symmetric matrix which satisfies the transition rule $h_{\alpha}=f_{\alpha \beta} h_{\beta} \bar{f}_{\alpha \beta}$. We define a connection 1-form locally by:

$$
\omega_{\alpha}=\partial h_{\alpha} h_{\alpha}^{-1}
$$

and compute the transition rule:

$$
\omega_{\alpha}=\partial\left\{f_{\alpha \beta} h_{\beta} \bar{f}_{\alpha \beta}\right\} \bar{f}_{\alpha \beta}^{-1} h_{\beta}^{-1} f_{\alpha \beta}^{-1}
$$

Since $f_{\alpha \beta}$ is holomorphic, $\bar{\partial} f_{\alpha \beta}=0$. This implies $d f_{\alpha \beta}=\partial f_{\alpha \beta}$ and $\partial \bar{f}_{\alpha \beta}=$ 0 so that:

$$
\omega_{\alpha}=\partial f_{\alpha \beta} f_{\alpha \beta}^{-1}+f_{\alpha \beta} \partial h_{\beta} h_{\beta}^{-1} f_{\alpha \beta}^{-1}=d f_{\alpha \beta} f_{\alpha \beta}^{-1}+f_{\alpha \beta} \omega_{\beta} f_{\alpha \beta}^{-1}
$$

Since this is the transition rule for a connection, the $\left\{\omega_{\alpha}\right\}$ patch together to define a connection $\nabla_{h}$.

It is immediate from the definition that:

$$
\left(\nabla_{h} s_{\alpha}, s_{\alpha}\right)+\left(s_{\alpha}, \nabla_{h} s_{\alpha}\right)=\omega_{\alpha} h_{\alpha}+h_{\alpha} \omega_{\alpha}^{*}=\partial h_{\alpha}+\bar{\partial} h_{\alpha}=d h_{\alpha}
$$

so $\nabla_{h}$ is a unitary connection on $V$. Since $\nabla_{h} s_{\alpha} \in C^{\infty}\left(\Lambda^{1,0} \otimes V\right)$ we conclude $\nabla_{h}$ vanishes on holomorphic sections when differentiated in antiholomorphic directions (i.e., $\nabla_{h}$ is a holomorphic connection). It is easily verified that these two properties determine $\nabla_{h}$.

We shall be particularly interested in holomorphic line bundles. If $L$ is a line bundle, then $h_{\alpha}$ is a positive function on $U_{\alpha}$ with $h_{\alpha}=\left|f_{\alpha \beta}\right|^{2} h_{\beta}$. The curvature in this case is a 2 -form defined by:

$$
\Omega_{\alpha}=d \omega_{\alpha}-\omega_{\alpha} \wedge \omega_{\alpha}=d\left(\partial h_{\alpha} h_{\alpha}^{-1}\right)=\bar{\partial} \partial \log \left(h_{\alpha}\right)=-\partial \bar{\partial} \log \left(h_{\alpha}\right)
$$

Therefore:

$$
c_{1}(L)=\frac{1}{2 \pi i} \partial \bar{\partial} \log (h)
$$

is independent of the holomorphic frame chosen for evaluation. $c_{1}(L) \in$ $C^{\infty}\left(\Lambda^{1,1}(M)\right)$ and $d c_{1}=\partial c_{1}=\bar{\partial} c_{1}=0$ so $c_{1}(L)$ is closed in all possible senses.

Let $\mathbf{C} P_{n}$ be the set of all lines through 0 in $\mathbf{C}^{n+1}$. Let $\mathbf{C}^{*}=\mathbf{C}-0$ act on $\mathbf{C}^{n+1}-0$ by complex multiplication, then $\mathbf{C} P_{n}=\left(\mathbf{C}^{n+1}-0\right) / \mathbf{C}^{*}$ and we give $\mathbf{C} P_{n}$ the quotient topology. Let $L$ be the tautological line bundle over $\mathbf{C} P_{n}$ :

$$
L=\left\{(x, z) \in \mathbf{C} P_{n} \times \mathbf{C}^{n+1}: z \in x\right\}
$$

Let $L^{*}$ be the dual bundle; this is called the hyperplane bundle.
Let $\left(z_{0}, \ldots, z_{n}\right)$ be the usual coordinates on $\mathbf{C}^{n+1}$. Let $U_{j}=\left\{z: z_{j} \neq\right.$ $0\}$. Since $U_{j}$ is $\mathbf{C}^{*}$ invariant, it projects to define an open set on $\mathbf{C} P_{n}$ we again shall denote by $U_{j}$. We define

$$
z_{k}^{j}=z_{k} / z_{j}
$$

over $U_{j}$. Since these functions are invariant under the action of $\mathbf{C}^{*}$, they extend to continuous functions on $U_{j}$. We let $z^{j}=\left(z_{0}^{j}, \ldots, \widehat{z_{j}^{j}}, \ldots, z_{n}^{j}\right)$ where we have deleted $z_{j}^{j}=1$. This gives local coordinates on $U_{j}$ in $\mathbf{C} P_{n}$. The transition relations are:

$$
z_{\nu}^{j}=\left(z_{j}^{k}\right)^{-1} z_{\nu}^{k}
$$

which are holomorphic so $\mathbf{C} P_{n}$ is a holomorphic manifold.
We let $s^{j}=\left(z_{0}^{j}, \ldots, z_{n}^{j}\right)$ be a section to $L$ over $U_{j}$. Then $s^{j}=\left(z_{j}^{k}\right)^{-1} s^{k}$ transform holomorphically so $L$ is a holomorphic line bundle over $\mathbf{C} P_{n}$. The coordinates $z_{j}$ on $\mathbf{C}^{n+1}$ give linear functions on $L$ and represent global holomorphic sections to the dual bundle $L^{*}$. There is a natural inner product on the trivial bundle $\mathbf{C} P_{n} \times \mathbf{C}^{n+1}$ which defines a fiber metric on $L$. We define:

$$
x=-c_{1}(L)=\frac{i}{2 \pi} \partial \bar{\partial} \log \left(1+\left|z^{j}\right|^{2}\right)
$$

then this is a closed 2 -form over $\mathbf{C} P_{n}$.
$\mathrm{U}(n+1)$ acts on $\mathbf{C}^{n+1}$; this action induces a natural action on $\mathbf{C} P_{n}$ and on $L$. Since the metric on $L$ arises from the invariant metric on $\mathbf{C}^{n+1}$, the 2-form $x$ is invariant under the action of $\mathrm{U}(n+1)$.
Lemma 2.3.1. Let $x=-c_{1}(L)$ over $\mathbf{C} P_{n}$. Then:
(a) $\int_{\mathbf{C} P_{n}} x^{n}=1$.
(b) $H^{*}\left(\mathbf{C} P_{n} ; \mathbf{C}\right)$ is a polynomial ring with additive generators $\left\{1, x, \ldots, x^{n}\right\}$.
(c) If $i: \mathbf{C} P_{n-1} \rightarrow \mathbf{C} P_{n}$ is the natural inclusion map, then $i^{*}(x)=x$.

Proof: (c) is immediate from the naturality of the constructions involved. Standard methods of algebraic topology give the additive structure of $H^{*}\left(\mathbf{C} P_{n} ; \mathbf{C}\right)=\mathbf{C} \oplus 0 \oplus \mathbf{C} \oplus 0 \cdots \oplus \mathbf{C}$. Since $x \in H^{2}\left(\mathbf{C} P_{n} ; \mathbf{C}\right)$ satisfies $x^{n} \neq 0$, (a) will complete the proof of (b). We fix a coordinate chart $U_{n}$. Since $\mathbf{C} P_{n}-U_{n}=\mathbf{C} P_{n-1}$, it has measure zero. It suffices to check that:

$$
\int_{U_{n}} x^{n}=1
$$

We identify $z \in U_{n}=\mathbf{C}^{n}$ with $(z, 1) \in \mathbf{C}^{n+1}$. We imbed $\mathrm{U}(n)$ in $\mathrm{U}(n+1)$ as the isotropy group of the vector $(0, \ldots, 0,1)$. Then $\mathrm{U}(n)$ acts on $(z, 1)$ in $\mathbf{C}^{n+1}$ exactly the same way as $\mathrm{U}(n)$ acts on $z$ in $\mathbf{C}^{n}$. Let dvol $=d x_{1} \wedge d y_{1} \wedge \cdots \wedge d x_{n} \wedge d y_{n}$ be ordinary Lebesgue measure on $\mathbf{C}^{n}$ without any normalizing constants of $2 \pi$. We parametrize $\mathbf{C}^{n}$ in the form $(r, \theta)$ for $0 \leq r<\theta$ and $\theta \in S^{2 n-1}$. We express $x^{n}=f(r, \theta)$ dvol. Since $x$ is invariant under the action of $\mathrm{U}(n), f(r, \theta)=f(r)$ is spherically symmetric and does not depend on the parameter $\theta$.

We compute that:

$$
\begin{aligned}
x & =-\frac{1}{2 \pi i} \partial \bar{\partial} \log \left(1+\sum z_{j} \bar{z}_{j}\right)=-\frac{1}{2 \pi i} \partial\left(\sum_{j} z_{j} d \bar{z}_{j} /\left(1+|z|^{2}\right)\right) \\
& =-\frac{1}{2 \pi i}\left\{\sum_{j} d z_{j} \wedge d \bar{z}_{j} /\left(1+r^{2}\right)-\sum_{j, k} z_{j} \bar{z}_{k} d z_{k} \wedge d \bar{z}_{j} /\left(1+r^{2}\right)^{2}\right\} .
\end{aligned}
$$

We evaluate at the point $z=(r, 0, \ldots, 0)$ to compute:

$$
\begin{aligned}
x & =-\frac{1}{2 \pi i}\left\{\sum_{j} d z_{j} d \bar{z}_{j} /\left(1+r^{2}\right)-r^{2} d z_{1} \wedge d \bar{z}_{1} /\left(1+r^{2}\right)^{2}\right\} \\
& =-\frac{1}{2 \pi i}\left\{d z_{1} \wedge d \bar{z}_{1} /\left(1+r^{2}\right)^{2}+\sum_{j>1} d z_{j} \wedge d \bar{z}_{j} /\left(1+r^{2}\right)\right\}
\end{aligned}
$$

Consequently at this point,

$$
\begin{aligned}
x^{n} & =\left(-\frac{1}{2 \pi i}\right)^{n} n!\left(1+r^{2}\right)^{-n-1} d z_{1} \wedge d \bar{z}_{1} \wedge \cdots \wedge d z_{n} \wedge d \bar{z}_{n} \\
& =\pi^{-n} n!\left(1+r^{2}\right)^{-n-1} d x_{1} \wedge d y_{1} \wedge \cdots \wedge d x_{n} \wedge d y_{n} \\
& =\pi^{-n} n!\left(1+r^{2}\right)^{-n-1} \text { dvol } .
\end{aligned}
$$

We use the spherical symmetry to conclude this identitiy holds for all $z$.
We integrate over $U_{n}=\mathbf{C}^{n}$ and use spherical coordinates:

$$
\begin{aligned}
\int x^{n} & =n!\pi^{-n} \int\left(1+r^{2}\right)^{-n-1} \mathrm{dvol}=n!\pi^{-n} \int\left(1+r^{2}\right)^{-n-1} r^{2 n-1} d r d \theta \\
& =n!\pi^{-n} \operatorname{vol}\left(S^{2 n-1}\right) \int_{0}^{\infty}\left(1+r^{2}\right)^{-n-1} r^{2 n-1} d r \\
& =\frac{1}{2} n!\pi^{-n} \operatorname{vol}\left(S^{2 n-1}\right) \int_{0}^{\infty}(1+t)^{-n-1} t^{n-1} d t
\end{aligned}
$$

We compute the volume of $S^{2 n-1}$ using the identity $\sqrt{\pi}=\int_{-\infty}^{\infty} e^{-t^{2}} d t$. Thus:

$$
\begin{aligned}
\pi^{n} & =\int e^{-r^{2}} \mathrm{dvol}=\operatorname{vol}\left(S^{2 n-1}\right) \int_{0}^{\infty} r^{2 n-1} e^{-r^{2}} d r \\
& =\frac{1}{2} \operatorname{vol}\left(S^{2 n-1}\right) \int_{0}^{\infty} t^{n-1} e^{-t} d t=\frac{(n-1)!}{2} \operatorname{vol}\left(S^{2 n-1}\right)
\end{aligned}
$$

We solve this for $\operatorname{vol}\left(S^{2 n-1}\right)$ and substitute to compute:

$$
\begin{aligned}
\int x^{n} & =n \int_{0}^{\infty}(1+t)^{-n-1} t^{n-1} d t=(n-1) \int_{0}^{\infty}(1+t)^{-n} t^{n-2} d t=\cdots \\
& =\int_{0}^{\infty}(1+t)^{-2} d t=1
\end{aligned}
$$

which completes the proof.
$x$ can also be used to define a $\mathrm{U}(n+1)$ invariant metric on $\mathbf{C} P_{n}$ called the Fubini-Study metric. We shall discuss this in more detail in Chapter 3 as this gives a Kaehler metric for $\mathbf{C} P_{n}$.

There is a relation between $L$ and $\Lambda^{1,0}\left(\mathbf{C} P_{n}\right)$ which it will be convenient to exploit:

Lemma 2.3.2. Let $M=\mathbf{C} P_{n}$.
(a) There is a short exact sequence of holomorphic bundles:

$$
0 \rightarrow \Lambda^{1,0}(M) \rightarrow L \otimes 1^{n+1} \rightarrow 1 \rightarrow 0
$$

(b) There is a natural isomorphism of complex bundles:

$$
T_{c}(M) \oplus 1 \simeq L^{*} \otimes 1^{n+1}=\underbrace{L^{*} \oplus \cdots \oplus L^{*}}_{n+1 \text { times }} .
$$

Proof: Rather than attempting to give a geometric proof of this fact, we give a combinatorial argument. Over $U_{k}$ we have functions $z_{i}^{k}$ which give coordinates $0 \leq i \leq n$ for $i \neq k$. Furthermore, we have a section $s^{k}$ to $L$. We let $\left\{s_{i}^{k}\right\}_{i=0}^{n}$ give a frame for $L \otimes 1^{n+1}=L \oplus \cdots \oplus L$. We define $F: \Lambda^{1,0}(M) \rightarrow L \otimes 1^{n+1}$ on $U_{k}$ by:

$$
F\left(d z_{i}^{k}\right)=s_{i}^{k}-z_{i}^{k} s_{k}^{k}
$$

We note that $z_{k}^{k}=1$ and $F\left(d z_{k}^{k}\right)=0$ so this is well defined on $U_{k}$. On the overlap, we have the relations $z_{i}^{k}=\left(z_{k}^{j}\right)^{-1} z_{i}^{j}$ and:

$$
s_{i}^{k}=\left(z_{k}^{j}\right)^{-1} s_{i}^{j} \quad \text { and } \quad d z_{i}^{k}=\left(z_{k}^{j}\right)^{-1} d z_{i}^{j}-\left(z_{k}^{j}\right)^{-2} z_{i}^{j} d z_{k}^{j} .
$$

Thus if we compute in the coordinate system $U_{j}$ we have:

$$
\begin{aligned}
F\left(d z_{i}^{k}\right) & =\left(z_{k}^{j}\right)^{-1}\left(s_{i}^{j}-z_{i}^{j} s_{j}^{j}\right)-\left(z_{k}^{j}\right)^{-2} z_{i}^{j}\left(s_{k}^{j}-z_{k}^{j} s_{j}^{j}\right) \\
& =s_{i}^{k}-z_{i}^{j} s_{j}^{k}-z_{i}^{k} s_{k}^{k}+z_{i}^{j} s_{j}^{k}=s_{i}^{k}-z_{i}^{k} s_{k}^{k}
\end{aligned}
$$

which agrees with the definition of $F$ on the coordinate system $U_{k}$. Thus $F$ is invariantly defined. It is clear $F$ is holomorphic and injective. We let
$\nu=L \otimes 1^{n+1} / \operatorname{image}(F)$, and let $\pi$ be the natural projection. It is clear that $s_{k}^{k}$ is never in the image of $F$ so $\pi s_{k}^{k} \neq 0$. Since

$$
F\left(d z_{j}^{k}\right)=s_{j}^{k}-z_{j}^{k} s_{k}^{k}=z_{j}^{k}\left(s_{j}^{j}-s_{k}^{k}\right)
$$

we conclude that $\pi s_{j}^{j}=\pi s_{k}^{k}$ so $s=\pi s_{k}^{k}$ is a globally defined non-zero section to $\nu$. This completes the proof of (a). We dualize to get a short exact sequence:

$$
0 \rightarrow 1 \rightarrow L^{*} \otimes 1^{n+1} \rightarrow T_{c}(M) \rightarrow 0
$$

These three bundles have natural fiber metrics. Any short exact sequence of vector bundles splits (although the splitting is not holomorphic) and this proves (b).

From assertion (b) it follows that the Chern class of $T_{c}\left(\mathbf{C} P_{n}\right)$ is given by:

$$
c\left(T_{c}\right)=c\left(T_{c} \oplus 1\right)=c\left(L^{*} \oplus \cdots \oplus L^{*}\right)=c\left(L^{*}\right)^{n+1}=(1+x)^{n+1}
$$

For example, if $m=1$, then $c\left(T_{c}\right)=(1+x)^{2}=1+2 x$. In this case, $c_{1}\left(T_{c}\right)=e\left(T\left(S^{2}\right)\right)$ and we computed $\int_{S^{2}} e(T(S))=2$ so $\int_{S^{2}} x=1$ which checks with Lemma 2.3.1.

If we forget the complex structure on $T_{c}(M)$ when $M$ is holomorphic, then we obtain the real tangent space $T(M)$. Consequently:

$$
T(M) \otimes \mathbf{C}=T_{c} \oplus T_{c}^{*}
$$

and

$$
\begin{aligned}
c\left(T\left(\mathbf{C} P_{n}\right) \otimes \mathbf{C}\right) & =c\left(T_{c}\left(\mathbf{C} P_{n}\right) \oplus T_{c}^{*}\left(\mathbf{C} P_{n}\right)\right) \\
& =(1+x)^{n+1}(1-x)^{n+1}=\left(1-x^{2}\right)^{n+1}
\end{aligned}
$$

When we take into account the sign changes involved in defining the total Pontrjagin form, we conclude:

Lemma 2.3.3.
(a) If $M=S^{2} \times \cdots \times S^{2}$ has dimension $2 n$, then $\int_{M} e(T(M))=\chi(M)=2^{n}$.
(b) If $M=\mathbf{C} P_{n}$ has dimension $2 n$, then

$$
c\left(T_{c}(M)\right)=(1+x)^{n+1} \quad \text { and } \quad p(T(M))=\left(1+x^{2}\right)^{n+1}
$$

$x \in H^{2}\left(\mathbf{C} P_{n} ; \mathbf{C}\right)$ is the generator given by $x=c_{1}\left(L^{*}\right)=-c_{1}(L), L$ is the tautological line bundle over $\mathbf{C} P_{n}$, and $L^{*}$ is the dual, the hyperplane bundle.

The projective spaces form a dual basis to both the real and complex characteristic classes. Let $\rho$ be a partition of the positive integer $k$ in the
form $k=i_{1}+\cdots+i_{j}$ where we choose the notation so $i_{1} \geq i_{2} \geq \cdots$. We let $\pi(k)$ denote the number of such partitions. For example, if $k=4$ then $\pi(4)=5$ and the possible partitions are:

$$
4=4, \quad 4=3+1, \quad 4=2+2, \quad 4=2+1+1, \quad 4=1+1+1+1
$$

We define classifying manifolds:

$$
M_{\rho}^{c}=\mathbf{C} P_{i_{1}} \times \cdots \times \mathbf{C} P_{i_{j}} \quad \text { and } \quad M_{\rho}^{r}=\mathbf{C} P_{2 i_{1}} \times \cdots \times \mathbf{C} P_{2 i_{j}}
$$

to be real manifolds of dimension $2 k$ and $4 k$.
Lemma 2.3.4. Let $k$ be a positive integer. Then:
(a) Let constants $c(\rho)$ be given. There exists a unique polynomial $P(A)$ of degree $k$ in the components of a $k \times k$ complex matrix which is GL $(k, \mathbf{C})$ invariant such that the characteristic class defined by $P$ satisfies:

$$
\int_{M_{\rho}^{c}} P\left(T_{c}\left(M_{\rho}^{c}\right)\right)=c(\rho)
$$

for every such partition $\rho$ of $k$.
(b) Let constants $c(\rho)$ be given. There exists a polynomial $P(A)$ of degree $2 k$ in the components of a $2 k \times 2 k$ real matrix which is $\mathrm{GL}(2 k, \mathbf{R})$ invariant such that the characteristic class defined by $P$ satisfies:

$$
\int_{M_{\rho}^{r}} P\left(T\left(M_{\rho}^{r}\right)\right)=c(\rho)
$$

for every such partition $\rho$ of $k$. If $P^{\prime}$ is another such polynomial, then $P(A)=P^{\prime}(A)$ for every skew-symmetric matrix $A$.

In other words, the real and complex characteristic classes are completely determined by their values on the appropriate classifying manifolds.

Proof: We prove (a) first. Let $\mathcal{P}_{k}$ denote the set of all such polynomials $P(A)$. We define $c_{\rho}=c_{i_{1}} \ldots c_{i_{j}} \in \mathcal{P}_{k}$, then by Lemma 2.1.3 the $\left\{c_{\rho}\right\}$ form a basis for $\mathcal{P}_{k}$ so $\operatorname{dim}\left(\mathcal{P}_{k}\right)=\pi(k)$. The $\left\{c_{\rho}\right\}$ are not a very convenient basis to work with. We will define instead:

$$
H_{\rho}=c h_{i_{1}} \ldots c h_{i_{j}} \in \mathcal{P}_{k}
$$

and show that the matrix:

$$
a(\rho, \tau)=\int_{M_{\tau}^{c}} H_{\rho}\left(T_{c}\left(M_{\tau}^{c}\right)\right)
$$

is a non-singular matrix. This will prove the $H_{\rho}$ also form a basis for $\mathcal{P}_{k}$ and that the $M_{\tau}^{c}$ are a dual basis. This will complete the proof.

The advantage of working with the Chern character rather than with the Chern class is that:

$$
c h_{i}\left(T_{c}\left(M_{1} \times M_{2}\right)\right)=c h_{i}\left(T_{c} M_{1} \oplus T_{c} M_{2}\right)=c h_{i}\left(T_{c} M_{1}\right)+c h_{i}\left(T_{c} M_{2}\right)
$$

Furthermore, $c h_{i}\left(T_{c} M\right)=0$ if $2 i>\operatorname{dim}(M)$. We define the length $\ell(\rho)=j$ to be the number of elements in the partition $\rho$. Then the above remarks imply:

$$
a(\rho, \tau)=0 \quad \text { if } \ell(\tau)>\ell(\rho)
$$

Furthermore, if $\ell(\tau)=\ell(\rho)$, then $a(\rho, \tau)=0$ unless $\tau=\rho$. We define the partial order $\tau>\rho$ if $\ell(\tau)>\ell(\rho)$ and extend this to a total order. Then $a(\rho, \tau)$ is a triangular matrix. To show it is invertible, it suffices to show the diagonal elements are non-zero.

We first consider the case in which $\rho=\tau=k$. Using the identity $T_{c}\left(\mathbf{C} P_{k}\right) \oplus 1=\left(L^{*} \oplus \cdots \oplus L^{*}\right)(k+1$ times $)$, it is clear that $c h_{k}\left(T_{c} \mathbf{C} P_{k}\right)=$ $(k+1) c h_{k}\left(L^{*}\right)$. For a line bundle, $c h_{k}\left(L^{*}\right)=c_{1}\left(L^{*}\right)^{k} / k!$. If $x=c_{1}\left(L^{*}\right)$ is the generator of $H^{2}\left(\mathbf{C} P_{k} ; \mathbf{C}\right)$, then $c h_{k}\left(T_{c} \mathbf{C} P_{k}\right)=(k+1) x^{k} / k$ ! which does not integrate to zero. If $\rho=\tau=\left\{i_{1}, \ldots, i_{j}\right\}$ then:

$$
a(\rho, \rho)=c \prod_{\nu=1}^{\ell(\rho)} a\left(i_{\nu}, i_{\nu}\right)
$$

where $c$ is a positive constant related to the multiplicity with which the $i_{\nu}$ appear. This completes the proof of (a).

To prove (b) we replace $M_{\rho}^{c}$ by $M_{2 \rho}^{c}=M_{\rho}^{r}$ and $H_{\rho}$ by $H_{2 \rho}$. We compute

$$
\begin{aligned}
c h_{2 i}(T(M)) \stackrel{\text { def }}{=} c h_{2 i}(T(M) \otimes \mathbf{C}) & =c h_{2 i}\left(T_{c} M \oplus T_{c}^{*} M\right) \\
& =c h_{2 i}\left(T_{c} M\right)+c h_{2 i}\left(T_{c}^{*} M\right) \\
& =2 c h_{2 i}\left(T_{c} M\right)
\end{aligned}
$$

Using this fact, the remainder of the proof of (b) is immediate from the calculations performed in (a) and this completes the proof.

The Todd class and the Hirzebruch $L$-polynomial were defined using generating functions. The generating functions were chosen so that they would be particularly simple on the classifying examples:

Lemma 2.3.5.
(a) Let $x_{j}=-\lambda_{j} / 2 \pi i$ be the normalized eigenvalues of a complex matrix A. We define $\operatorname{Td}(A)=\operatorname{Td}(x)=\prod_{j} x_{j} /\left(1-e^{-x_{j}}\right)$ as the Todd class. Then:

$$
\int_{M_{\rho}^{c}} T d\left(T_{c}\left(M_{\rho}^{c}\right)\right)=1 \quad \text { for all } \rho
$$

(b) Let $x_{j}=\lambda_{j} / 2 \pi$ where the eigenvalues of the skew-symmetric real ma$\operatorname{trix} A$ are $\left\{ \pm \lambda_{1}, \ldots\right\}$. We define $L(A)=L(x)=\prod_{j} x_{j} / \tanh x_{j}$ as the Hirzebruch L-polynomial. Then $\int_{M_{\rho}^{r}} L\left(T\left(M^{r}\right)\right)=1$ for all $\rho$.
We will use this calculation to prove the integral of the Todd class gives the arithmetic genus of a complex manifold and that the integral of the Hirzebruch $L$-polynomial gives the signature of an oriented real manifold. In each case, we only integrate the part of the total class which is of the same degree as the dimension of the manifold.

Proof: $T d$ is a multiplicative class:

$$
\operatorname{Td}\left(T_{c}\left(M_{1} \times M_{2}\right)\right)=\operatorname{Td}\left(T_{c}\left(M_{1}\right) \oplus T_{c}\left(M_{2}\right)\right)=\operatorname{Td}\left(T_{c}\left(M_{1}\right)\right) \operatorname{Td}\left(T_{c}\left(M_{2}\right)\right)
$$

Similarly the Hirzebruch polynomial is multiplicative. This shows it suffices to prove Lemma 2.3.5 in the case $\rho=k$ so $M=\mathbf{C} P_{k}$ or $\mathbf{C} P_{2 k}$.

We use the decomposition $T_{c}\left(\mathbf{C} P_{k}\right) \oplus 1=L^{*} \oplus \cdots \oplus L^{*} \quad(k+1$ times $)$ to compute $T d\left(T_{c}\left(\mathbf{C} P_{k}\right)\right)=[T d(x)]^{k+1}$ where $x=c_{1}\left(L^{*}\right)$ is the generator of $H^{2}\left(\mathbf{C} P_{k} ; \mathbf{R}\right)$. Since $x^{k}$ integrates to 1 , it suffices to show the coefficient of $x^{k}$ in $\operatorname{Td}(x)^{k+1}$ is 1 or equivalently to show:

$$
\operatorname{Res}_{x=0} x^{-k-1}\left[T d(x)^{k+1}\right]=\operatorname{Res}_{x=0}\left(1-e^{-x}\right)^{-k-1}=1 .
$$

If $k=0$, then:

$$
\left(1-e^{-x}\right)^{-1}=\left(x-\frac{1}{2} x^{2}+\cdots\right)^{-1}=x^{-1}\left(1+\frac{1}{2} x+\cdots\right)
$$

and the result follows. Similarly, if $k=1$

$$
\left(1-e^{-x}\right)^{-2}=x^{-2}(1+x+\cdots)
$$

and the result follows. For larger values of $k$, proving this directly would be a combinatorial nightmare so we use instead a standard trick from complex variables. If $g(x)$ is any meromorphic function, then $\operatorname{Res}_{x=0} g^{\prime}(x)=0$. We apply this to the function $g(x)=\left(1-e^{-x}\right)^{-k}$ for $k \geq 1$ to conclude:

$$
\operatorname{Res}_{x=0}\left(1-e^{-x}\right)^{-k-1} e^{-x}=0
$$

This implies immediately that:

$$
\begin{aligned}
\operatorname{Res}_{x=0}\left(1-e^{-x}\right)^{-k-1} & =\operatorname{Res}_{x=0}\left(1-e^{-x}\right)^{-k-1}\left(1-e^{-x}\right) \\
& =\operatorname{Res}_{x=0}\left(1-e^{-x}\right)^{-k}=1
\end{aligned}
$$

by induction which completes the proof of assertion (a).

We now assume $k$ is even and study $\mathbf{C} P_{k}$. Again, using the decomposition $T_{c}\left(\mathbf{C} P_{k}\right) \oplus 1=L^{*} \oplus \cdots \oplus L^{*}$ it follows that

$$
L\left(T\left(\mathbf{C} P_{k}\right)\right)=[L(x)]^{k+1}=\frac{x^{k+1}}{\tanh ^{k+1} x}
$$

Since we are interested in the coefficient of $x^{k}$, we must show:

$$
\operatorname{Res}_{x=0} \tanh ^{-k-1} x=1 \quad \text { if } k \text { is even }
$$

or equivalently

$$
\operatorname{Res}_{x=0} \tanh ^{-k} x=1 \quad \text { if } k \text { is odd. }
$$

We recall that

$$
\begin{aligned}
\tanh x & =\frac{e^{x}-e^{-x}}{e^{x}+e^{-x}} \\
\tanh ^{-1} x & =\frac{2+x^{2}+\cdots}{2 x+\cdots}
\end{aligned}
$$

so the result is clear if $k=1$. We now proceed by induction. We differentiate $\tanh ^{-k} x$ to compute:

$$
\left(\tanh ^{-k} x\right)^{\prime}=-k\left(\tanh ^{-k-1} x\right)\left(1-\tanh ^{2} x\right)
$$

This implies

$$
\operatorname{Res}_{x=0} \tanh ^{-k+1} x=\operatorname{Res}_{x=0} \tanh ^{-k-1} x
$$

for any integer $k$. Consequently $\operatorname{Res}_{x=0} \tanh ^{-k} x=1$ for any odd integer $k$ since these residues are periodic modulo 2. (The residue at $k$ even is, of course, zero). This completes the proof.

### 2.4. The Gauss-Bonnet Theorem.

Let $P$ be a self-adjoint elliptic partial differential operator of order $d>0$. If the leading symbol of $P$ is positive definite, we derived an asymptotic expansion for $\operatorname{Tr}\left\{e^{-t P}\right\}$ in section 1.7. This is too general a setting in which to work so we shall restrict attention henceforth to operators with leading symbol given by the metric tensor, as this is the natural category in which to work.

Let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be a second order operator. We choose a local frame for $V$. Let $x=\left(x_{1}, \ldots, x_{m}\right)$ be a system of local coordinates. Let $d s^{2}=g_{i j} d x^{i} d x^{j}$ be the metric tensor and let $g^{i j}$ denote the inverse matrix; $d x^{i} \cdot d x^{j}=g^{i j}$ is the metric on the dual space $T^{*}(M)$. We assume $P$ has the form:

$$
P=-g^{i j} \frac{\partial^{2}}{\partial x_{i} \partial x_{j}} I+a_{j} \frac{\partial}{\partial x_{j}}+b
$$

where we sum over repeated indices. The $a_{j}$ and $b$ are sections to $\operatorname{END}(V)$.
We introduce formal variables for the derivatives of the symbol of $P$. Let

$$
g_{i j / \alpha}=d_{x}^{\alpha} g_{i j}, \quad a_{j / \alpha}=d_{x}^{\alpha} a_{j}, \quad b_{/ \alpha}=d_{x}^{\alpha} b
$$

We will also use the notation $g_{i j / k l \ldots,}, a_{j / k l \ldots}$, and $b_{/ k l \ldots}$... for multiple partial derivatives. We emphasize that these variables are not tensorial but depend upon the choice of a coordinate system (and local frame for $V$ ).

There is a natural grading on these variables. We define:

$$
\operatorname{ord}\left(g_{i j / \alpha}\right)=|\alpha|, \quad \operatorname{ord}\left(a_{j / \alpha}\right)=1+|\alpha|, \quad \operatorname{ord}\left(b_{/ \alpha}\right)=2+|\alpha|
$$

Let $\mathcal{P}$ be the non-commutative polynomial algebra in the variables of positive order. We always normalize the coordinate system so $x_{0}$ corresponds to the point $(0, \ldots, 0)$ and so that:

$$
g_{i j}\left(x_{0}\right)=\delta_{i j} \quad \text { and } \quad g_{i j / k}\left(x_{0}\right)=0
$$

Let $K(t, x, x)$ be the kernel of $e^{-t P}$. Expand

$$
K(t, x, x) \sim \sum_{n \geq 0} t^{\frac{n-m}{2}} e_{n}(x, P)
$$

where $e_{n}(x, P)$ is given by Lemma 1.7.4. Lemma 1.7.5(c) implies:
Lemma 2.4.1. $e_{n}(x, P) \in \mathcal{P}$ is a non-commutative polynomial in the jets of the symbol of $P$ which is homogeneous of order $n$. If $a_{n}(x, P)=$ $\operatorname{Tr} e_{n}(x, P)$ is the fiber trace, then $a_{n}(x, P)$ is homogeneous of order $n$ in
the components (relative to some local frame) of the jets of the total symbol of $P$.

We could also have proved Lemma 2.4.1 using dimensional analysis and the local nature of the invariants $e_{n}(x, P)$ instead of using the combinatorial argument given in Chapter 1.

We specialize to the case where $P=\Delta_{p}$ is the Laplacian on $p$-forms. If $p=0$, then $\Delta_{0}=-g^{-1} \partial / \partial x_{i}\left\{g g^{i j} \partial / \partial x_{j}\right\} ; \quad g=\operatorname{det}\left(g_{i j}\right)^{1 / 2}$ defines the Riemannian volume dvol $=g d x$. The leading symbol is given by the metric tensor; the first order symbol is linear in the 1-jets of the metric with coefficients which depend smoothly on the metric tensor. The $0^{\text {th }}$ order symbol is zero in this case. More generally:
Lemma 2.4.2. Let $x=\left(x_{1}, \ldots, x_{m}\right)$ be a system of local coordinates on $M$. We use the $d x^{I}$ to provide a local frame for $\Lambda\left(T^{*} M\right)$. Relative to this frame, we expand $\sigma\left(\Delta_{p}\right)=p_{2}+p_{1}+p_{0} . \quad p_{2}=|\xi|^{2} I . \quad p_{1}$ is linear in the 1-jets of the metric with coefficients which depend smoothly on the $g_{i j}$ 's. $p_{0}$ is the sum of a term which is linear in the 2-jets of the metric and a term which is quadratic in the 1-jets of the metric with coefficients which depend smoothly on the $g_{i j}$ 's.
Proof: We computed the leading symbol in the first chapter. The remainder of the lemma follows from the decomposition $\Delta=d \delta+\delta d=$ $\pm d * d * \pm * d * d$. In flat space, $\Delta=-\sum_{i} \partial^{2} / \partial x_{i}^{2}$ as computed earlier. If the metric is curved, we must also differentiate the matrix representing the Hodge * operator. Each derivative applied to "*" reduces the order of differentiation by one and increases the order in the jets of the metric by one.

Let $a_{n}(x, P)=\operatorname{Tr} e_{n}(x, P)$. Then:

$$
\operatorname{Tr} e^{-t P} \sim \sum_{n \geq 0} t^{\frac{n-m}{2}} \int_{M} a_{n}(x, P) \operatorname{dvol}(x)
$$

For purposes of illustration, we give without proof the first few terms in the asymptotic expansion of the Laplacian. We shall discuss such formulas in more detail in the fourth chapter.
Lemma 2.4.3.
(a) $a_{0}\left(x, \Delta_{p}\right)=(4 \pi)^{-1} \operatorname{dim}\left(\Lambda^{p}\right)=(4 \pi)^{-1}\binom{m}{p}$.
(b) $a_{2}\left(x, \Delta_{0}\right)=(4 \pi)^{-1}\left(-R_{i j i j}\right) / 6$.
(c)
$a_{4}\left(x, \Delta_{0}\right)=\frac{1}{4 \pi} \cdot \frac{-12 R_{i j i j ; k k}+5 R_{i j i j} R_{k l k l}-2 R_{i j i k} R_{l j l k}+2 R_{i j k l} R_{i j k l}}{360}$.
In this lemma, ";" denotes multiple covariant differentiation. $a_{0}, a_{2}$, $a_{4}$, and $a_{6}$ have been computed for $\Delta_{p}$ for all $p$, but as the formulas are extremely complicated, we shall not reproduce them here.

We consider the algebra generated by the variables $\left\{g_{i j / \alpha}\right\}_{|\alpha| \geq 2}$. If $X$ is a coordinate system and $G$ a metric and if $P$ is a polynomial in these variables, we define $P(X, G)\left(x_{0}\right)$ by evaluation. We always normalize the choice of $X$ so $g_{i j}(X, G)\left(x_{0}\right)=\delta_{i j}$ and $g_{i j / k}(X, G)\left(x_{0}\right)=0$, and we omit these variables from consideration. We say that $P$ is invariant if $P(X, G)\left(x_{0}\right)=P(Y, G)\left(x_{0}\right)$ for any two normalized coordinate systems $X$ and $Y$. We denote the common value by $P(G)\left(x_{0}\right)$. For example, the scalar curvature $K=-\frac{1}{2} R_{i j i j}$ is invariant as are the $a_{n}\left(x, \Delta_{p}\right)$.

We let $\mathcal{P}_{m}$ denote the ring of all invariant polynomials in the derivatives of the metric for a manifold of dimension $m$. We defined $\operatorname{ord}\left(g_{i j / \alpha}\right)=|\alpha|$; let $\mathcal{P}_{m, n}$ be the subspace of invariant polynomials which are homogeneous of order $n$. This is an algebraic characterization; it is also useful to have the following coordinate free characterization:
Lemma 2.4.4. Let $P \in \mathcal{P}_{m}$, then $P \in \mathcal{P}_{m, n}$ if and only if $P\left(c^{2} G\right)\left(x_{0}\right)=$ $c^{-n} P(G)\left(x_{0}\right)$ for every $c \neq 0$.
Proof: Fix $c=0$ and let $X$ be a normalized coordinate system for the metric $G$ at the point $x_{0}$. We assume $x_{0}=(0, \ldots, 0)$ is the center of the coordinate system $X$. Let $Y=c X$ be a new coordinate system, then:

$$
\begin{aligned}
\partial / \partial y_{i} & =c^{-1} \partial / \partial x_{i} & c^{2} G\left(\partial / \partial y_{i}, \partial / \partial y_{j}\right) & =G\left(\partial / \partial x_{i}, \partial / \partial x_{j}\right) \\
d_{y}^{\alpha} & =c^{-|\alpha|} d_{x}^{\alpha} & g_{i j / \alpha}\left(Y, c^{2} G\right) & =c^{-|\alpha|} g_{i j / \alpha}(X, G) .
\end{aligned}
$$

This implies that if $A$ is any monomial of $P$ that:

$$
A\left(Y, c^{2} G\right)\left(x_{0}\right)=c^{-\operatorname{ord}(A)} A(X, G)\left(x_{0}\right)
$$

Since $Y$ is normalized coordinate system for the metric $c^{2} G, P\left(c^{2} G\right)\left(x_{0}\right)=$ $P\left(Y, c^{2} G\right)\left(x_{0}\right)$ and $P(G)\left(x_{0}\right)=P(X, G)\left(x_{0}\right)$. This proves the Lemma.

If $P \in \mathcal{P}_{m}$ we can always decompose $P=P_{0}+\cdots+P_{n}$ into homogeneous polynomials. Lemma 2.4.4 implies the $P_{j}$ are all invariant separately. Therefore $\mathcal{P}_{m}$ has a direct sum decomposition $\mathcal{P}_{m}=\mathcal{P}_{m, 0} \oplus \mathcal{P}_{m, 1} \oplus \cdots \oplus$ $\mathcal{P}_{m, n} \oplus \cdots$ and has the structure of a graded algebra. Using Taylor's theorem, we can always find a metric with the $g_{i j / \alpha}(X, G)\left(x_{0}\right)=c_{i j, \alpha}$ arbitrary constants for $|\alpha| \geq 2$ and so that $g_{i j}(X, G)\left(x_{0}\right)=\delta_{i j}, g_{i j / k}(X, G)\left(x_{0}\right)=0$. Consequently, if $P \in \mathcal{P}_{m}$ is non-zero as a polynomial, then we can always find $G$ so $P(G)\left(x_{0}\right) \neq 0$ so $P$ is non-zero as a formula. It is for this reason we work with the algebra of jets. This is a pure polynomial algebra. If we work instead with the algebra of covariant derivatives of the curvature tensor, we must introduce additional relations which correspond to the Bianchi identities as this algebra is not a pure polynomial algebra.

We note finally that $\mathcal{P}_{m, n}$ is zero if $n$ is odd since we may take $c=-1$. Later in this chapter, we will let $\mathcal{P}_{m, n, p}$ be the space of $p$-form valued invariants which are homogeneous of order $n$. A similar argument will show $\mathcal{P}_{m, n, p}$ is zero if $n-p$ is odd.

Lemmas 2.4.1 and 2.4.2 imply:

Lemma 2.4.5. $a_{n}\left(x, \Delta_{p}\right)$ defines an element of $\mathcal{P}_{m, n}$.
This is such an important fact that we give another proof based on Lemma 2.4.4 to illustrate the power of dimensional analysis embodied in this lemma. Fix $c>0$ and let $\Delta_{p}\left(c^{2} G\right)=c^{-2} \Delta_{p}(G)$ be the Laplacian corresponding to the new metric. Since $\operatorname{dvol}\left(c^{2} G\right)=c^{m} \operatorname{dvol}(G)$, we conclude:

$$
\begin{aligned}
e^{-t \Delta_{p}\left(c^{2} G\right)} & =e^{-t c^{-2} \Delta_{p}(G)} \\
K\left(t, x, x, \Delta_{p}\left(c^{2} G\right)\right) \mathrm{dvol}\left(c^{2} G\right) & =K\left(c^{-2} t, x, x, \Delta_{p}(G)\right) \operatorname{dvol}(G) \\
K\left(t, x, x, \Delta_{p}\left(c^{2} G\right)\right) & =c^{-m} K\left(c^{-2} t, x, x, \Delta_{p}(G)\right) \\
\sum_{n} t^{\frac{n-m}{2}} a_{n}\left(x, \Delta_{p}\left(c^{2} G\right)\right) & \sim \sum_{n} c^{-m} c^{m} c^{-n} t^{\frac{n-m}{2}} a_{n}\left(x, \Delta_{p}(G)\right) \\
a_{n}\left(x, \Delta_{p}\left(c^{2} G\right)\right) & =c^{-n} a_{n}\left(x, \Delta_{p}(G)\right)
\end{aligned}
$$

We expand $a_{n}\left(x, \Delta_{p}(G)\right)=\sum_{\nu} a_{n, \nu, p}+r$ in a finite Taylor series about $g_{i j}=\delta_{i j}$ in the $g_{i j / \alpha}$ variables. Then if $a_{n, \nu, p}$ is the portion which is homogeneous of order $\nu$, we use this identity to show $a_{n, \nu, p}=0$ for $n \neq \nu$ and to show the remainder in the Tayor series is zero. This shows $a_{n}$ is a homogeneous polynomial of order $n$ and completes the proof.

Since $a_{n}\left(x, \Delta_{p}\right)=0$ if $n$ is odd, in many references the authors replace the asymptotic series by $t^{-\frac{m}{2}} \sum_{n} t^{n} a_{n}\left(x, \Delta_{p}\right)$, They renumber this sequence $a_{0}, a_{1}, \ldots$ rather than $a_{0}, 0, a_{2}, 0, a_{4}, \ldots$ We shall not adopt this notational convention as it makes dealing with boundary problems more cumbersome.
H. Weyl's theorem on the invariants of the orthogonal group gives a spanning set for the spaces $\mathcal{P}_{m, n}$ :
Lemma 2.4.6. We introduce formal variables $R_{i_{1} i_{2} i_{3} i_{4} ; i_{5} \ldots i_{k}}$ for the multiple covariant derivatives of the curvature tensor. The order of such a variable is $k+2$. We consider the polynomial algebra in these variables and contract on pairs of indices. Then all possible such expressions generate $\mathcal{P}_{m}$. In particular:

$$
\begin{gathered}
\{1\} \text { spans } \mathcal{P}_{m, 0}, \quad\left\{R_{i j i j}\right\} \text { spans } \mathcal{P}_{m, 2} \\
\left\{R_{i j i j ; k k}, R_{i j i j} R_{k l k l}, R_{i j i k} R_{l j l k}, R_{i j k l} R_{i j k l}\right\} \text { spans } \mathcal{P}_{m, 4}
\end{gathered}
$$

This particular spanning set for $\mathcal{P}_{m, 4}$ is linearly independent and forms a basis if $m \geq 4$. If $m=3, \operatorname{dim}\left(\mathcal{P}_{3,4}\right)=3$ while if $m=2, \operatorname{dim}\left(\mathcal{P}_{2,4}\right)=2$ so there are relations imposed if the dimension is low. The study of these additional relations is closely related to the Gauss-Bonnet theorem.

There is a natural restriction map

$$
r: \mathcal{P}_{m, n} \rightarrow \mathcal{P}_{m-1, n}
$$

which is defined algebraically as follows. We let

$$
\operatorname{deg}_{k}\left(g_{i j / \alpha}\right)=\delta_{i, k}+\delta_{j, k}+\alpha(k)
$$

be the number of times an index $k$ appears in the variable $g_{i j / \alpha}$. Let

$$
r\left(g_{i j / \alpha}\right)= \begin{cases}g_{i j / \alpha} \in \mathcal{P}_{m-1} & \text { if } \operatorname{deg}_{m}\left(g_{i j / \alpha}\right)=0 \\ 0 & \text { if } \operatorname{deg}_{m}\left(g_{i j / \alpha}\right)>0\end{cases}
$$

We extend $r: \mathcal{P}_{m} \rightarrow \mathcal{P}_{m-1}$ to be an algebra morphism; $r(P)$ is a polynomial in the derivatives of a metric on a manifold of dimension $m-1$. It is clear $r$ preserves the degree of homogeneity.
$r$ is the dual of a natural extension map. Let $G^{\prime}$ be a metric on a manifold $M^{\prime}$ of dimension $m-1$. We define $i\left(G^{\prime}\right)=G+d \theta^{2}$ on the manifold $M=M^{\prime} \times S^{1}$ where $S^{1}$ is the unit circle with natural parameter $\theta$. If $X^{\prime}$ is a local coordinate system on $M^{\prime}$, then $i\left(X^{\prime}\right)=(x, \theta)$ is a local coordinate system on $M$. It is clear that:

$$
(r P)\left(X^{\prime}, G^{\prime}\right)\left(x_{0}^{\prime}\right)=P\left(i\left(X^{\prime}\right), i\left(G^{\prime}\right)\right)\left(x_{0}^{\prime}, \theta_{0}\right)
$$

for any $\theta_{0} \in S^{1}$; what we have done by restricting to product manifolds $M^{\prime} \times S^{1}$ with product metrics $G^{\prime}+d \theta^{2}$ is to introduce the relation which says the metric is flat in the last coordinate. Restiction is simply the dual of this natural extension; $r P$ is invariant if $P$ is invariant. Therefore $r: \mathcal{P}_{m, n} \rightarrow \mathcal{P}_{m, n-1}$.

There is one final description of the restriction map which will be useful. In discussing a H . Weyl spanning set, the indices range from 1 through $m$. We define the restriction by letting the indices range from 1 through $m-1$. Thus $R_{i j i j} \in \mathcal{P}_{m, 2}$ is its own restriction in a formal sense; of course $r\left(R_{i j i j}\right)=0$ if $m=2$ since there are no non-trivial local invariants over a circle.

Theorem 2.4.7.
(a) $r: \mathcal{P}_{m, n} \rightarrow \mathcal{P}_{m-1, n}$ is always surjective.
(b) $r: \mathcal{P}_{m, n} \rightarrow \mathcal{P}_{m-1, n}$ is bijective if $n<m$.
(c) $r: \mathcal{P}_{m, m} \rightarrow \mathcal{P}_{m-1, m}$ has 1-dimensional kernel spanned by the Euler class $E_{m}$ if $m$ is even. If $m$ is odd, $\mathcal{P}_{m, m}=\mathcal{P}_{m-1, m}=0$.
(c) is an axiomatic characterization of the Euler form. It is an expression of the fact that the Euler form is a unstable characteristic class as opposed to the Pontrjagin forms which are stable characteristic classes.
Proof: (a) is consequence of H . Weyl's theorem. If we choose a H . Weyl spanning set, we let the indices range from 1 to $m$ instead of from 1 to $m-1$ to construct an element in the inverse image of $r$. The proof of (b) and of (c) is more complicated and will be postponed until the next subsection. Theorem 2.4.7 is properly a theorem in invariance theory, but we have stated it at this time to illustrate how invariance theory can be used to prove index theorems using heat equation methods:

ThEOREM 2.4.8. Let $a_{n}(x, d+\delta)=\sum_{p}(-1)^{p} a_{n}\left(x, \Delta_{p}\right) \in \mathcal{P}_{n, m}$ be the invariant of the de Rham complex. We showed in Lemma 1.7.6 that:

$$
\int_{M} a_{n}(x, d+\delta) \operatorname{dvol}(x)= \begin{cases}\chi(M) & \text { if } n=m \\ 0 & \text { if } n \neq m\end{cases}
$$

Then:
(a) $a_{n}(x, d+\delta)=0$ if either $m$ is odd or if $n<m$.
(b) $a_{m}(x, d+\delta)=E_{m}$ is $m$ is even so $\chi(M)=\int_{M} E_{m} \operatorname{dvol}(x)$ (Gauss-

## Bonnet theorem).

Proof: We suppose first that $m$ is odd. Locally we can always choose an orientation for $T(M)$. We let $*$ be the Hodge operator then $* \Delta_{p}= \pm \Delta_{m-p} *$ locally. Since these operators are locally isomorphic, their local invariants are equal so $a_{n}\left(x, \Delta_{p}\right)=a_{n}\left(x, \Delta_{m-p}\right)$. If $m$ is odd, these terms cancel in the alternating sum to give $a_{n}(x, d+\delta)=0$. Next we suppose $m$ is even. Let $M=M^{\prime} \times S^{1}$ with the product metric. We decompose any $\omega \in \Lambda\left(T^{*} M\right)$ uniquely in the form

$$
\omega=\omega_{1}+\omega_{2} \wedge d \theta \quad \text { for } \omega_{i} \in \Lambda\left(T^{*} M^{\prime}\right)
$$

We define:

$$
F(\omega)=\omega_{1} \wedge d \theta+\omega_{2}
$$

and compute easily that $F \Delta=\Delta F$ since the metric is flat in the $S^{1}$ direction. If we decompose $\Lambda(M)=\Lambda^{\mathrm{e}}(M) \oplus \Lambda^{\circ}(M)$ into the forms of even and odd degree, then $F$ interchanges these two factors. Therefore, $a_{n}\left(x, \Delta_{e}\right)=a_{n}\left(x, \Delta_{0}\right)$ so $a_{n}(x, d+\delta)=0$ for such a product metric. This implies $r\left(a_{n}\right)=0$. Therefore $a_{n}=0$ for $n<m$ by Theorem 2.4.7. Furthermore:

$$
a_{m}=c_{m} E_{m}
$$

for some universal constant $c_{m}$. We show $c_{m}=1$ by integrating over the classifying manifold $M=S^{2} \times \cdots \times S^{2}$. Let $2 \bar{m}=m$, then

$$
2^{\bar{m}}=\chi(M)=\int_{M} a_{m}(x, d+\delta) \operatorname{dvol}(x)=\int_{M} E_{m} \operatorname{dvol}(x)
$$

by Lemma 2.3.4. This completes the proof of the Gauss-Bonnet theorem.

### 2.5. Invariance Theory and the Pontrjagin Classes of the Tangent Bundle.

In the previous subsection, we gave in Theorem 2.4.7 an axiomatic characterization of the Euler class in terms of functorial properties. In this subsection we will complete the proof of Theorem 2.4.7. We will also give a similar axiomatic characterization of the Pontrjagin classes which we will use in our discussion of the signature complex in Chapter 3.

Let $T: \mathbf{R}^{m} \rightarrow \mathbf{R}^{m}$ be the germ of a diffeomorphism. We assume that:

$$
T(0)=0 \quad \text { and } \quad d T(x)=d T(0)+O\left(x^{2}\right) \quad \text { for } d T(0) \in \mathrm{O}(k)
$$

If $X$ is any normalized coordinate system for a metric $G$, then $T X$ is another normalized coordinate system for $G$. We define an action of the group of germs of diffeomorphisms on the polynomial algebra in the $\left\{g_{i j / \alpha}\right\},|\alpha| \geq 2$ variables by defining the evaluation:

$$
\left(T^{*} P\right)(X, G)\left(x_{0}\right)=P(T X, G)\left(x_{0}\right)
$$

Clearly $P$ is invariant if and only if $T^{*} P=P$ for every such diffeomorphism $T$.

Let $P$ be invariant and let $A$ be a monomial. We let $c(A, P)$ be the coefficient of $A$ in $P ; \quad c(A, P)$ defines a linear functional on $\mathcal{P}$ for any monomial $A$. We say $A$ is a monomial of $P$ if $c(A, P) \neq 0$. Let $T_{j}$ be the linear transformation:

$$
T_{j}\left(x_{k}\right)= \begin{cases}-x_{j} & \text { if } k=j \\ x_{k} & \text { if } k \neq j\end{cases}
$$

This is reflection in the hyperplane defined by $x_{j}=0$. Then

$$
T_{j}^{*}(A)=(-1)^{\operatorname{deg}_{j}(A)} A
$$

for any monomial $A$. Since

$$
T_{j}^{*} P=\sum(-1)^{\operatorname{deg}_{j}(A)} c(A, P) A=P=\sum c(A, P) A
$$

we conclude $\operatorname{deg}_{j}(A)$ must be even for any monomial $A$ of $P$. If $A$ has the form:

$$
A=g_{i_{1} j_{1} / \alpha_{1}} \ldots g_{i_{r} j_{r} / \alpha_{r}}
$$

we define the length of $A$ to be:

$$
\ell(A)=r .
$$

It is clear $2 \ell(A)+\operatorname{ord}(A)=\sum_{j} \operatorname{deg}_{j}(A)$ so $\operatorname{ord}(A)$ is necessarily even if $A$ is a monomial of $P$. This provides another proof $\mathcal{P}_{m, n}=0$ if $n$ is odd.

In addition to the hyperplane reflections, it is convenient to consider coordinate permutations. If $\rho$ is a permutation, then $T_{\rho}^{*}(A)=A^{\rho}$ is defined by replacing each index $i$ by $\rho(i)$ in the variables $g_{i j / k \ldots .}$. Since $T_{\rho}^{*}(P)=P$, we conclude the form of $P$ is invariant under coordinate permutations or equivalently $c(A, P)=c\left(A^{\rho}, P\right)$ for every monomial $A$ of $P$.

We can use these two remarks to begin the proof of Theorem 2.4.7. Fix $P \neq 0$ with $r(P)=0$ and $P \in \mathcal{P}_{m, n}$. Let $A$ be a monomial of $P$. Then $r(P)=0$ implies $\operatorname{deg}_{m}(A)>0$ is even. Since $P$ is invariant under coordinate permutations, $\operatorname{deg}_{k}(A) \geq 2$ for $1 \leq k \leq m$. We construct the chain of inequalities:

$$
\begin{aligned}
2 m \leq \sum_{k} \operatorname{deg}_{k}(A) & =2 \ell(A)+\operatorname{ord}(A) \\
& =\sum_{\nu}\left(2+\left|\alpha_{\nu}\right|\right) \leq \sum_{\nu} 2\left|\alpha_{\nu}\right|=2 \operatorname{ord}(A)=2 n
\end{aligned}
$$

We have used the fact $\left|\alpha_{\nu}\right| \geq 2$ for $1 \leq \nu \leq \ell(A)$. This implies $m \leq n$ so in particular $P=0$ if $n<m$ which proves assertion (b) of Theorem 2.4.7. If $n=m$, then all of the inequalities in this string must be equalities. This implies $2 \ell(A)=m, \operatorname{deg}_{k}(A)=2$ for all $k$, and $\left|\alpha_{\nu}\right|=2$ for all $\nu$. Thus the monomial $A$ must have the form:

$$
A=g_{i_{1} j_{1} / k_{1} l_{1}} \ldots g_{i_{r} j_{r} / k_{r} l_{r}} \quad \text { for } 2 r=m
$$

and in particular, $P$ only involves the second derivatives of the metric.
In order to complete the proof of Theorem 2.4.7, we must use some results involving invariance under circle actions:

Lemma 2.5.1. Parametrize the circle $\mathrm{SO}(2)$ by $z=(a, b)$ for $a^{2}+b^{2}=1$. Let $T_{z}$ be the coordinate transformation:

$$
y_{1}=a x_{1}+b x_{2}, \quad y_{2}=-b x_{1}+a x_{2}, \quad y_{k}=x_{k} \quad \text { for } k>2
$$

Let $P$ be a polynomial and assume $T_{z}^{*} P=P$ for all $z \in S^{1}$. Then:
(a) If $g_{12 / \alpha}$ divides a monomial $A$ of $P$ for some $\alpha$, then $g_{11 / \beta}$ divides a monomial $B$ of $P$ for some $\beta$.
(b) If $g_{i j / \alpha}$ divides a monomial of $A$ of $P$ for some $(i, j)$, then $g_{k l / \beta}$ divides a monomial $B$ of $P$ for some $\beta$ and some $(k, l)$ where $\beta(1)=\alpha(1)+\alpha(2)$ and $\beta(2)=0$.
Of course, the use of the indices 1 and 2 is for convenience only. This lemma holds true for any pair of indices under the appropriate invariance assumption.

We postpone the proof of this lemma for the moment and use it to complete the proof of Theorem 2.4.7. Let $P \neq 0 \in \mathcal{P}_{m, m}$ with $r(P)=0$.

Let $A$ be a monomial of $P$. We noted $\operatorname{deg}_{k}(A)=2$ for $1 \leq k \leq m$ and $A$ is a polynomial in the 2 -jets of the metric. We decompose

$$
A=g_{i_{1} j_{1} / k_{1} l_{1}} \ldots g_{i_{r} j_{r} / k_{r} l_{r}} \quad \text { for } 2 r=m
$$

By making a coordinate permutation we can choose $A$ so $i_{1}=1$. If $j_{1} \neq 1$, we can make a coordinate permutation to assume $j_{1}=2$ and then apply Lemma 2.5.1(a) to assume $i_{1}=j_{1}=1$. We let $P_{1}=\sum c(A, P) A$ where the sum ranges over $A$ of the form:

$$
A=g_{11 / k_{1} l_{1}} \ldots g_{i_{r} j_{r} / k_{r} l_{r}}
$$

which are monomials of $P$. Then $P_{1} \neq 0$ and $P_{1}$ is invariant under coordinate transformations which fix the first coordinate. Since $\operatorname{deg}_{1}(A)=2$, the index 1 appears nowhere else in $A$. Thus $k_{1}$ is not 1 so we may make a coordinate permutation to choose $A$ so $k_{1}=2$. If $l_{2} \neq 2$, then $l_{1} \geq 3$ so we may make a coordinate permutation to assume $l_{1}=3$. We then apply Lemma 2.5.1(b) to choose $A$ a monomial of $P_{1}$ of the form:

$$
A=g_{11 / 22} A^{\prime}
$$

We have $\operatorname{deg}_{1}(A)=\operatorname{deg}_{2}(A)=2$ so $\operatorname{deg}_{1}\left(A^{\prime}\right)=\operatorname{deg}_{2}\left(A^{\prime}\right)=0$ so these indices do not appear in $A^{\prime}$. We define $P_{2}=\sum c(A, P) A$ where the sum ranges over those monomial $A$ of $P$ which are divisible by $g_{11 / 22}$, then $P_{2} \neq 0$.

We proceed inductively in this fashion to show finally that

$$
A_{0}=g_{11 / 22} g_{33 / 44} \ldots g_{m-1, m-1 / m m}
$$

is a monomial of $P$ so $c\left(A_{0}, P\right) \neq 0$. The function $c\left(A_{0}, P\right)$ is a separating linear functional on the kernel of $r$ in $\mathcal{P}_{m, m}$ and therefore

$$
\operatorname{dim}\left(\left\{P \in \mathcal{P}_{m, m}: r(P)=0\right\}\right) \leq 1
$$

We complete the proof of Theorem 2.4.7(c) by showing that $r\left(E_{m}\right)=0$. If $E_{m}$ is the Euler form and $M=M^{\prime} \times S^{1}$, then:

$$
E_{m}(M)=E_{m}\left(M^{\prime} \times S^{1}\right)=E_{m-1}\left(M^{\prime}\right) E_{1}\left(S^{1}\right)=0
$$

which completes the proof; $\operatorname{dim} \mathrm{N}(r)=1$ and $E_{m} \in \mathrm{~N}(r)$ spans.
Before we begin the proof of Lemma 2.5.1, it is helpful to consider a few examples. If we take $A=g_{11 / 11}$ then it is immediate:

$$
\begin{aligned}
& \partial / \partial y_{1}=a \partial / \partial x_{1}+b \partial / \partial x_{2}, \quad \partial / \partial y_{2}=-b \partial / \partial x_{1}+a \partial / \partial x_{2} \\
& \partial / \partial y_{k}=\partial / \partial x_{k} \quad \text { for } k>2 .
\end{aligned}
$$

We compute $T_{z}^{*}(A)$ by formally replacing each index 1 by $a 1+b 2$ and each index 2 by $-b 1+a 2$ and expanding out the resulting expression. Thus, for example:

$$
\begin{aligned}
T_{z}^{*}\left(g_{11 / 11}\right)= & g_{a 1+b 2, a 1+b 2 / a 1+b 2, a 1+b 2} \\
= & a^{4} g_{11 / 11}+2 a^{3} b g_{11 / 12}+2 a^{3} b g_{12 / 11} \\
& \quad+a^{2} b^{2} g_{11 / 22}+a^{2} b^{2} g_{22 / 11}+4 a^{2} b^{2} g_{12 / 12} \\
& \quad+2 a b^{3} g_{12 / 22}+2 a b^{3} g_{22 / 12}+b^{4} g_{22 / 22}
\end{aligned}
$$

We note that those terms involving $a^{3} b$ arose from changing exactly one index to another index $(1 \rightarrow 2$ or $2 \rightarrow 1)$; the coefficient reflects the multiplicity. Thus in particular this polynomial is not invariant.

In computing invariance under the circle action, all other indices remain fixed so

$$
\begin{aligned}
T_{z}^{*}\left(g_{34 / 11}+g_{34 / 22}\right)= & \left(a^{2} g_{34 / 11}+b^{2} g_{34 / 22}+2 a b g_{34 / 12}\right. \\
& \left.\quad+a^{2} g_{34 / 22}+b^{2} g_{34 / 11}-2 a b g_{34 / 12}\right) \\
= & \left(a^{2}+b^{2}\right)\left(g_{34 / 11}+g_{34 / 22}\right)
\end{aligned}
$$

is invariant. Similarly, it is easy to compute:

$$
T_{z}^{*}\left(g_{11 / 22}+g_{22 / 11}-2 g_{12 / 12}\right)=\left(a^{2}+b^{2}\right)^{2}\left(g_{11 / 22}+g_{22 / 11}-2 g_{12 / 12}\right)
$$

so this is invariant. We note this second example is homogeneous of degree 4 in the $(a, b)$ variables since $\operatorname{deg}_{1}(A)+\operatorname{deg}_{2}(A)=4$.

With these examples in mind, we begin the proof of Lemma 2.5.1. Let $P$ be invariant under the action of the circle acting on the first two coordinates. We decompose $P=P_{0}+P_{1}+\cdots$ where each monomial $A$ of $P_{j}$ satisfies $\operatorname{deg}_{1}(A)+\operatorname{deg}_{2}(A)=j$. If $A$ is such a monomial, then $T_{z}^{*}(A)$ is a sum of similar monomials. Therefore each of the $P_{j}$ is invariant separately so we may assume $P=P_{n}$ for some $n$. By setting $a=b=-1$, we see $n$ must be even. Decompose:

$$
T_{z}^{*}(P)=a^{n} P^{(0)}+b a^{n-1} P^{(1)}+\cdots+b^{n} P^{(n)}
$$

where $P=P^{(0)}$. We use the assumption $T_{z}^{*}(P)=P$ and replace $b$ by $-b$ to see:

$$
0=T_{(a, b)}^{*}(P)-T_{(a,-b)}^{*}(P)=2 b a^{n-1} P^{(1)}+2 b^{3} a^{n-3} P^{(3)}+\cdots
$$

We divide this equation by $b$ and take the limit as $b \rightarrow 0$ to show $P^{(1)}=0$. (In fact, it is easy to show $P^{(2 j+1)}=0$ and $P^{(2 j)}=\binom{n / 2}{j} P$ but as we shall not need this fact, we omit the proof).

We let $A_{0}$ denote a variable monomial which will play the same role as the generic constant $C$ of Chapter 1 . We introduce additional notation we shall find useful. If $A$ is a monomial, we decompose $T_{(a, b)}^{*} A=a^{n} A+$ $a^{n-1} b A^{(1)}+\cdots$. If $B$ is a monomial of $A^{(1)}$ then $\operatorname{deg}_{1}(B)=\operatorname{deg}_{1}(A) \pm 1$ and $B$ can be constructed from the monomial $A$ by changing exactly one index $1 \rightarrow 2$ or $2 \rightarrow 1$. If $\operatorname{deg}_{1}(B)=\operatorname{deg}_{1}(A)+1$, then $B$ is obtained from $A$ by changing exactly one index $2 \rightarrow 1 ; c\left(B, A^{(1)}\right)$ is a negative integer which reflects the multiplicity with which this change can be made. If $\operatorname{deg}_{1}(B)=\operatorname{deg}_{1}(A)-1$, then $B$ is obtained from $A$ by changing exactly one index $1 \rightarrow 2 ; c\left(B, A^{(1)}\right)$ is a positive integer. We define:

$$
\begin{array}{rlrl}
A(1 \rightarrow 2) & =\sum c\left(B, A^{(1)}\right) B & \text { summed over } \operatorname{deg}_{1}(B)=\operatorname{deg}_{1}(A)-1 \\
A(2 \rightarrow 1) & =\sum-c\left(B, A^{(1)}\right) B & \text { summed over } \operatorname{deg}_{1}(B)=\operatorname{deg}_{1}(A)+1 \\
A^{(1)} & =A(1 \rightarrow 2)-A(2 \rightarrow 1) .
\end{array}
$$

For example, if $A=\left(g_{12 / 33}\right)^{2} g_{11 / 44}$ then $n=6$ and:

$$
\begin{aligned}
& A(1 \rightarrow 2)=2 g_{12 / 33} g_{22 / 33} g_{11 / 44}+2\left(g_{12 / 33}\right)^{2} g_{12 / 44} \\
& A(2 \rightarrow 1)=2 g_{12 / 33} g_{11 / 33} g_{11 / 44}
\end{aligned}
$$

It is immediate from the definition that:

$$
\sum_{B} c(B, A(1 \rightarrow 2))=\operatorname{deg}_{1}(A) \quad \text { and } \quad \sum_{B} c(B, A(2 \rightarrow 1))=\operatorname{deg}_{2}(A)
$$

Finally, it is clear that $c\left(B, A^{(1)}\right) \neq 0$ if and only if $c\left(A, B^{(1)}\right) \neq 0$, and that these two coefficients will be opposite in sign, and not necessarily equal in magnitude.

Lemma 2.5.2. Let $P$ be invariant under the action of $\mathrm{SO}(2)$ on the first two coordinates and let $A$ be a monomial of $P$. Let $B$ be a monomial of $A^{(1)}$. Then there exists a monomial of $A_{1}$ different from $A$ so that

$$
c\left(B, A^{(1)}\right) c(A, P) c\left(B, A_{1}^{(1)}\right) c\left(A_{1}, P\right)<0 .
$$

Proof: We know $P^{(1)}=0$. We decompose

$$
P^{(1)}=\sum_{A} c(A, P) A^{(1)}=\sum_{B} c(A, P) c\left(B, A^{(1)}\right) B
$$

Therefore $c\left(B, P^{(1)}\right)=0$ implies

$$
\sum_{A} c(A, P) c\left(B, A^{(1)}\right)=0
$$

for all monomials $B$. If we choose $B$ so $c\left(B, A^{(1)}\right) \neq 0$ then there must be some other monomial $A_{1}$ of $P$ which helps to cancel this contribution. The signs must be opposite which proves the lemma.

This lemma is somewhat technical and formidable looking, but it is exactly what we need about orthogonal invariance. We can now complete the proof of Lemma 2.5.1. Suppose first that $A=\left(g_{12 / \alpha}\right)^{k} A_{0}$ for $k>0$ where $g_{12 / \alpha}$ does not divide $A_{0}$. Assume $c(A, P) \neq 0$. Let $B=g_{11 / \alpha}\left(g_{12 / \alpha}\right)^{k-1} A_{0}$ then $c\left(B, A^{(1)}\right)=-k \neq 0$. Choose $A_{1} \neq A$ so $c\left(A_{1}, P\right) \neq 0$. Then $c\left(A_{1}, B^{(1)}\right) \neq 0$. If $c\left(A_{1}, B(2 \rightarrow 1)\right) \neq 0$ then $A_{1}$ is constructed from $B$ by changing a $2 \rightarrow 1$ index so $A_{1}=g_{11 / \beta} \ldots$ has the desired form. If $c\left(A_{1}, B(1 \rightarrow 2)\right) \neq 0$, we expand $B(1 \rightarrow 2)=A+$ terms divisible by $g_{11 / \beta}$ for some $\beta$. Since $A_{1} \neq A$, again $A_{1}$ has the desired form which proves (a).

To prove (b), we choose $g_{i j / \alpha}$ dividing some monomial of $P$. Let $\beta$ be chosen with $\beta(1)+\beta(2)=\alpha(1)+\alpha(2)$ and $\beta(k)=\alpha(k)$ for $k>2$ with $\alpha(1)$ maximal so $g_{u v / \beta}$ divides some monomial of $P$ for some $(u, v)$. Suppose $\beta(2) \neq$ 0 , we argue for a contradiction. Set $\gamma=(\beta(1)+1, \beta(2)-1, \beta(3), \ldots, \beta(m))$. Expand $A=\left(g_{u v / \beta}\right)^{k} A_{0}$ and define $B=g_{u v / \gamma}\left(g_{u v / \beta}\right)^{k-1} A_{0}$ where $g_{u v / \beta}$ does not divide $A_{0}$. Then $c\left(B, A^{(1)}\right)=-\beta(2) k \neq 0$ so we may choose $A_{1} \neq A$ so $c\left(A_{1}, B^{(1)}\right) \neq 0$. If $c\left(A_{1}, B(2 \rightarrow 1)\right) \neq 0$ then either $A_{1}$ is divisible by $g_{u^{\prime} v^{\prime} / \gamma}$ or by $g_{u v / \gamma^{\prime}}$ where $\gamma^{\prime}(1)=\gamma(1)+1, \quad \gamma^{\prime}(2)=\gamma(2)-1$, and $\gamma^{\prime}(j)=\alpha(j)$ for $j>2$. Either possibility contradicts the choice of $\beta$ as maximal so $c\left(A_{1}, B(1 \rightarrow 2)\right) \neq 0$. However, $B(1 \rightarrow 2)=\beta(1) A+$ terms divisible by $g_{u^{\prime} v^{\prime} / \gamma}$ for some $\left(u^{\prime}, v^{\prime}\right)$. This again contradicts the maximality as $A \neq A_{1}$ and completes the proof of Lemma 2.5.1 and thereby of Theorems 2.4.7 and 2.4.8.

If $I=\left\{1 \leq i_{1} \leq \cdots \leq i_{p} \leq m\right\}$, let $|I|=p$ and $d x^{I}=d x_{i_{1}} \wedge \cdots \wedge$ $d x_{i_{p}}$. A $p$-form valued polynomial is a collection $\left\{P_{I}\right\}=P$ for $|I|=p$ of polynomials $P_{I}$ in the derivatives of the metric. We will also sometimes write $P=\sum_{|I|=p} P_{I} d x^{I}$ as a formal sum to represent $P$. If all the $\left\{P_{I}\right\}$ are homogeneous of order $n$, we say $P$ is homogeneous of order $n$. We define:

$$
P(X, G)\left(x_{0}\right)=\sum_{I} P_{I}(X, G) d x^{I} \in \Lambda^{p}\left(T^{*} M\right)
$$

to be the evaluation of such a polynomial. We say $P$ is invariant if $P(X, G)\left(x_{0}\right)=P(Y, G)\left(x_{0}\right)$ for every normalized coordinate systems $X$ and $Y$; as before we denote the common value by $P(G)\left(x_{0}\right)$. In analogy with Lemma 2.4.4 we have:
Lemma 2.5.3. Let $P$ be $p$-form valued and invariant. Then $P$ is homogeneous of order $n$ if and only if $P\left(c^{2} G\right)\left(x_{0}\right)=c^{p-n} P(G)\left(x_{0}\right)$ for every $c \neq 0$.

The proof is exactly the same as that given for Lemma 2.4.4. The only new feature is that $d y^{I}=c^{p} d x^{I}$ which contributes the extra feature of $c^{p}$ in this equation.

We define $\mathcal{P}_{m, n, p}$ to be the vector space of all $p$-form valued invariants which are homogeneous of order $n$ on a manifold of dimension $m$. If $\mathcal{P}_{m, *, p}$ denotes the vector space of all $p$-form valued invariant polynomials, then we have a direct sum decomposition $\mathcal{P}_{m, *, p}=\bigoplus_{n} \mathcal{P}_{m, n, p}$ exactly as in the scalar case $p=0$.

We define $\operatorname{deg}_{k}(I)$ to be 1 if $k$ appears in $I$ and 0 if $k$ does not appear in $I$.

Lemma 2.5.4. Let $P \in \mathcal{P}_{m, n, p}$ with $P \neq 0$. Then $n-p$ is even. If $A$ is a monomial then $A$ is a monomial of at most one $P_{I}$. If $c\left(A, P_{I}\right) \neq 0$ then:

$$
\operatorname{deg}_{k}(A)+\operatorname{deg}_{k}(I) \quad \text { is always even. }
$$

Proof: Let $T$ be the coordinate transformation defined by $T\left(x_{k}\right)=-x_{k}$ and $T\left(x_{j}\right)=x_{j}$ for $j \neq k$. Then

$$
P=T^{*}(P)=\sum_{I, A}(-1)^{\operatorname{deg}_{k}(A)+\operatorname{deg}_{k}(I)} c\left(A, P_{I}\right) d x^{I}
$$

which implies $\operatorname{deg}_{k}(A)+\operatorname{deg}_{k}(I)$ is even if $c\left(A, P_{I}\right) \neq 0$. Therefore

$$
\sum_{k} \operatorname{deg}_{k}(A)+\operatorname{deg}_{k}(I)=2 \ell(A)+\operatorname{ord}(A)+p=2 \ell(A)+n+p
$$

must be even. This shows $n+p$ is even if $P \neq 0$. Furthermore, if $c\left(A, P_{I}\right) \neq$ 0 then $I$ is simply the ordered collection of indices $k \operatorname{sog}_{k}(A)$ is odd which shows $A$ is a monomial of at most one $P_{I}$ and completes the proof.

We extend the Riemannian metric to a fiber metric on $\Lambda\left(T^{*} M\right) \otimes \mathbf{C}$. It is clear that $P \cdot P=\sum_{I} P_{I} \bar{P}_{I}$ since the $\left\{d x^{I}\right\}$ form an orthonormal basis at $x_{0} . \quad P \cdot P$ is a scalar invariant. We use Lemma 2.5.1 to prove:

Lemma 2.5.5. Let $P$ be $p$-form valued and invariant under the action of $\mathrm{O}(m)$. Let $A$ be a monomial of $P$. Then there is a monomial $A_{1}$ of $P$ with $\operatorname{deg}_{k}\left(A_{1}\right)=0$ for $k>2 \ell(A)$.

Proof: Let $r=\ell(A)$ and let $P_{r}^{\prime}=\sum_{\ell(B)=r} c\left(B, P_{I}\right) d x^{I} \neq 0$. Since this is invariant under the action of $\mathrm{O}(m)$, we may assume without loss of generality $P=P_{r}$. We construct a scalar invariant by taking the inner product $Q=(P, P)$. By applying Lemma 2.5.1(a) and making a coordinate permutation if necessary, we can assume $g_{11 / \alpha_{1}}$ divides some monomial of $Q$. We apply 2.5.1(b) to the indices $>1$ to assume $\alpha_{1}(k)=0$ for $k>2$. $g_{11 / \alpha_{1}}$ must divide some monomial of $P$. Let

$$
P_{1}=\sum_{A_{0}, I} c\left(g_{11 / \alpha_{1}} A_{0}, P_{I}\right) g_{11 / \alpha_{1}} A_{0} d x^{I} \neq 0
$$

This is invariant under the action of $\mathrm{O}(m-2)$ on the last $m-2$ coordinates. We let $Q_{1}=\left(P_{1}, P_{1}\right)$ and let $g_{i_{2} j_{2} / \alpha_{2}}$ divide some monomial of $Q_{1}$. If both $i_{2}$ and $j_{2}$ are $\leq 2$ we leave this alone. If $i_{2} \leq 2$ and $j_{2} \geq 3$ we perform a coordinate permutation to assume $j_{2}=3$. In a similar fashion if $i_{2} \geq 3$ and $j_{2} \leq 2$, we choose this variable so $i_{2}=3$. Finally, if both indices $\geq 3$, we apply Lemma 2.5.1(a) to choose this variable so $i_{2}=j_{2}=3$. We apply Lemma 2.5.1(b) to the variables $k \geq 4$ to choose this variable so $\alpha_{2}(k)=0$ for $k>4$. If $A_{2}=g_{11 / \alpha_{1}} g_{i_{2} j_{2} / \alpha_{2}}$ then:

$$
\operatorname{deg}_{k}\left(A_{2}\right)=0 \text { for } k>4 \text { and } A_{2} \text { divides some monomial of } P
$$

We continue inductively to construct $A_{r}=g_{11 / \alpha_{1}} \ldots g_{i_{r} j_{r} / \alpha_{r}}$ so that

$$
\operatorname{deg}_{k}\left(A_{r}\right)=0 \text { for } k>2 r \text { and } A_{r} \text { divides some monomial of } P .
$$

Since every monomial of $P$ has length $r$, this implies $A_{r}$ itself is a monomial of $P$ and completes the proof.

Let $P_{j}(G)=p_{j}(T M)$ be the $j^{\text {th }}$ Pontrjagin form computed relative to the curvature tensor of the Levi-Civita connection. If we expand $p_{j}$ in terms of the curvature tensor, then $p_{j}$ is homogeneous of order $2 j$ in the $\left\{R_{i j k l}\right\}$ tensor so $p_{j}$ is homogeneous of order $4 j$ in the jets of the metric. It is clear $p_{j}$ is an invariantly defined $4 j$-form so $p_{j} \in \mathcal{P}_{m, 4 j, 4 j}$. The algebra generated by the $p_{j}$ is called the algebra of the Pontrjagin forms. By Lemma 2.2.2, this is also the algebra of real characteristic forms of $T(M)$. If $\rho$ is a partition of $k=i_{1}+\cdots+i_{j}$ we define $p_{\rho}=p_{i_{1}} \ldots p_{i_{j}} \in \mathcal{P}_{m, 4 k, 4 k}$. The $\left\{p_{\rho}\right\}$ form a basis of the Pontrjagin $4 k$ forms. By Lemma 2.3.4, these are linearly independent if $m=4 k$ since the matrix $\int_{M_{\tau}^{r}} p_{\rho}$ is non-singular. By considering products of these manifolds with flat tori $T^{m-4 k}$ we can easily show that the $\left\{p_{\rho}\right\}$ are linearly independent in $\mathcal{P}_{m, 4 k, 4 k}$ if $4 k \leq m$. We let $\pi(k)$ be the number of partitions of $k$; this is the dimension of the Pontrjagin forms.

The axiomatic characterization of the real characteristic forms of the tangent space which is the analogue of the axiomatic characterization of the Euler class given in Theorem 2.4.7 is the following:

Lemma 2.5.6.
(a) $\mathcal{P}_{m, n, p}=0$ if $n<p$.
(b) $\mathcal{P}_{m, n, n}$ is spanned by the Pontrjagin forms-i.e.,

$$
\begin{gathered}
\mathcal{P}_{m, n, n}=0 \text { if } n \text { is not divisible by } 4 k \\
\mathcal{P}_{m, 4 k, 4 k}=\operatorname{span}\left\{p_{\rho}\right\} \text { for } 4 k \leq m \text { has dimension } \pi(k)
\end{gathered}
$$

Proof: By decomposing $P$ into its real and imaginary parts, it suffices to prove this lemma for polynomials with real coefficients. Let $0 \neq P \in \mathcal{P}_{m, n, p}$
and let $A$ be a monomial of $P$. Use Lemma 2.5.5 to find a monomial $A_{1}$ of some $P_{I}$ where $I=\left\{1 \leq i_{1}<\cdots<i_{p} \leq m\right\}$ so $\operatorname{deg}_{k}(A)=0$ for $k>2 \ell(A)$. Since $\operatorname{deg}_{i_{p}}(A)$ is odd (and hence non-zero) and since $2 \ell(A) \leq \sum\left|\alpha_{\nu}\right|=n$ as $A$ is a polynomial in the jets of order 2 and higher, we estimate:

$$
p \leq i_{p} \leq 2 \ell(A) \leq n
$$

This proves $P=0$ if $n<p$, which proves (a).
If $n=p$, then all of these inequalities must have been equalities. This shows that the higher order jets of the metric do not appear in $P$ so $P$ is a polynomial in the $\left\{g_{i j / k l}\right\}$ variables. Furthermore, $i_{p}=p$ and there is some monomial $A$ so that

$$
\operatorname{deg}_{k}(A)=0 \text { for } k>p \quad \text { and } \quad A d x^{1} \wedge \cdots \wedge d x^{p} \text { appears in } P .
$$

There is a natural restriction map $r: \mathcal{P}_{m, n, p} \rightarrow \mathcal{P}_{m-1, n, p}$ defined in the same way as the restriction map $r: \mathcal{P}_{m, n, 0} \rightarrow \mathcal{P}_{m-1, n, 0}$ discussed earlier. This argument shows $r: \mathcal{P}_{m, n, n} \rightarrow \mathcal{P}_{m-1, n, n}$ is injective for $n<m$ since $r\left(A d x^{1} \wedge \cdots \wedge d x^{p}\right)=A d x^{1} \wedge \cdots \wedge d x^{p}$ appears in $r(P)$. The Pontrjagin forms have dimension $\pi(k)$ for $n=4 k$. By induction, $r^{m-n}: \mathcal{P}_{m, n, n} \rightarrow$ $\mathcal{P}_{n, n, n}$ is injective so $\operatorname{dim} \mathcal{P}_{m, n, n} \leq \operatorname{dim} \mathcal{P}_{n, n, n}$.

We shall prove Lemma 2.5.6(b) for the special case $n=m$. If $n$ is not divisible by $4 k$, then $\operatorname{dim} \mathcal{P}_{n, n, n}=0$ which implies $\operatorname{dim} \mathcal{P}_{m, n, n}=0$. If $n=4 k$, then $\operatorname{dim} \mathcal{P}_{n, n, n}=\pi(k)$ implies that $\pi(k) \leq \operatorname{dim} \mathcal{P}_{m, n, n} \leq$ $\operatorname{dim} \mathcal{P}_{n, n, n} \leq \pi(k)$ so $\operatorname{dim} \mathcal{P}_{m, n, n}=\pi(k)$. Since the Pontrjagin forms span a subspace of exactly dimension $\pi(k)$ in $\mathcal{P}_{m, n, n}$ this will complete the proof of (b).

This lemma is at the heart of our discussion of the index theorem. We shall give two proofs for the case $n=m=p$. The first is based on H. Weyl's theorem for the orthogonal group and follows the basic lines of the proof given in Atiyah-Bott-Patodi. The second proof is purely combinatorial and follows the basic lines of the original proof first given in our thesis. The H. Weyl based proof has the advantage of being somewhat shorter but relies upon a deep theorem we have not proved here while the second proof although longer is entirely self-contained and has some additional features which are useful in other applications.

We review H. Weyl's theorem briefly. Let $V$ be a real vector space with a fixed inner product. Let $\mathrm{O}(V)$ denote the group of linear maps of $V \rightarrow V$ which preserve this inner product. Let $\otimes^{k}(V)=V \otimes \cdots \otimes V$ denote the $k^{\text {th }}$ tensor prodect of $V$. If $g \in \mathrm{O}(V)$, we extend $g$ to act orthogonally on $\bigotimes^{k}(V)$. We let $z \mapsto g(z)$ denote this action. Let $f: \otimes^{k}(V) \rightarrow \mathbf{R}$ be a multi-linear map, then we say $f$ is $\mathrm{O}(V)$ invariant if $f(g(z))=f(z)$ for every $g \in \mathrm{O}(V)$. By letting $g=-1$, it is easy to see there are no $\mathrm{O}(V)$ invariant maps if $k$ is odd. We let $k=2 j$ and construct a map
$f_{0}: \otimes^{k}(V)=(V \otimes V) \otimes(V \otimes V) \otimes \cdots \otimes(V \otimes V) \rightarrow \mathbf{R}$ using the metric to map $(V \otimes V) \rightarrow \mathbf{R}$. More generally, if $\rho$ is any permutation of the integers 1 through $k$, we define $z \mapsto z_{\rho}$ as a map from $\bigotimes^{k}(V) \rightarrow \bigotimes^{k}(V)$ and let $f_{\rho}(z)=f_{0}\left(z_{\rho}\right)$. This will be $\mathrm{O}(V)$ invariant for any permutation $\rho$. H. Weyl's theorem states that the maps $\left\{f_{\rho}\right\}$ define a spanning set for the collection of $\mathrm{O}(V)$ invariant maps.

For example, let $k=4$. Let $\left\{v_{i}\right\}$ be an orthonormal basis for $V$ and express any $z \in \bigotimes^{4}(V)$ in the form $a_{i j k l} v_{i} \otimes v_{j} \otimes v_{k} \otimes v_{l}$ summed over repeated indices. Then after weeding out duplications, the spanning set is given by:

$$
f_{0}(z)=a_{i i j j}, \quad f_{1}(z)=a_{i j i j}, \quad f_{2}(z)=a_{i j j i}
$$

where we sum over repeated indices. $f_{0}$ corresponds to the identity permutation; $f_{1}$ corresponds to the permutation which interchanges the second and third factors; $f_{2}$ corresponds to the permutation which interchanges the second and fourth factors. We note that these need not be linearly independent; if $\operatorname{dim} V=1$ then $\operatorname{dim}\left(\bigotimes^{4} V\right)=1$ and $f_{1}=f_{2}=f_{3}$. However, once $\operatorname{dim} V$ is large enough these become linearly independent.

We are interested in $p$-form valued invariants. We take $\bigotimes^{k}(V)$ where $k-p$ is even. Again, there is a natural map we denote by

$$
f^{p}(z)=f_{0}\left(z_{1}\right) \wedge \Lambda\left(z_{2}\right)
$$

where we decompose $\bigotimes^{k}(V)=\bigotimes^{k-p}(V) \otimes \bigotimes^{p}(V)$. We let $f_{0}$ act on the first $k-p$ factors and then use the natural map $\bigotimes^{p}(V) \xrightarrow{\Lambda} \Lambda^{p}(V)$ on the last $p$ factors. If $\rho$ is a permutation, we set $f_{\rho}^{p}(z)=f^{p}\left(z_{\rho}\right)$. These maps are equivariant in the sense that $f_{\rho}^{p}(g z)=g f_{\rho}^{p}(z)$ where we extend $g$ to act on $\Lambda^{p}(V)$ as well. Again, these are a spanning set for the space of equivariant multi-linear maps from $\bigotimes^{k}(V)$ to $\Lambda^{p}(V)$. If $k=4$ and $p=2$, then after eliminating duplications this spanning set becomes:

$$
\begin{array}{lll}
f_{1}(z)=a_{i i j k} v_{j} \wedge v_{k}, & f_{2}(z)=a_{i j i k} v_{j} \wedge v_{k}, & f_{3}(z)=a_{i j k i} v_{j} \wedge v_{k} \\
f_{4}(z)=a_{j i k i} v_{j} \wedge v_{k}, & f_{5}(z)=a_{j i i k} v_{i} \wedge v_{k}, & f_{6}(z)=a_{j k i i} v_{j} \wedge v_{k}
\end{array}
$$

Again, these are linearly independent if $\operatorname{dim} V$ is large, but there are relations if $\operatorname{dim} V$ is small. Generally speaking, to construct a map from $\bigotimes^{k}(V) \rightarrow \Lambda^{p}(V)$ we must alternate $p$ indices (the indices $j, k$ in this example) and contract the remaining indices in pairs (there is only one pair $i, i$ here).
Theorem 2.5.7 (H. Weyl's theorem on the invariants of the orthogonal group). The space of maps $\left\{f_{\rho}^{p}\right\}$ constructed above span the space of equivariant multi-linear maps from $\bigotimes^{k} V \rightarrow \Lambda^{p} V$.

The proof of this theorem is beyond the scope of the book and will be omitted. We shall use it to give a proof of Lemma 2.5.6 along the lines
of the proof given by Atiyah, Bott, and Patodi. We will then give an independent proof of Lemma 2.5.6 by other methods which does not rely on H. Weyl's theorem.

We apply H. Weyl's theorem to our situation as follows. Let $P \in \mathcal{P}_{n, n, n}$ then $P$ is a polynomial in the 2 -jets of the metric. If we let $X$ be a system of geodesic polar coordinates centered at $x_{0}$, then the 2 -jets of the the metric are expressible in terms of the curvature tensor so we can express $P$ as a polynomial in the $\left\{R_{i j k l}\right\}$ variables which is homogeneous of order $n / 2$. The curvature defines an element $R \in \bigotimes^{4}(T(M))$ since it has 4 indices. There are, however, relations among the curvature variables:

$$
R_{i j k l}=R_{k l i j}, \quad R_{i j k l}=-R_{j i k l}, \quad \text { and } \quad R_{i j k l}+R_{i k l j}+R_{i l j k}=0
$$

We let $V$ be the sub-bundle of $\bigotimes^{4}(T(M))$ consisting of tensors satisfying these 3 relations.

If $n=2$, then $P: V \rightarrow \Lambda^{2}(T(M))$ is equivariant while more generally, $P: \bigotimes^{n / 2}(V) \rightarrow \Lambda^{n}(T(M))$ is equivariant under the action of $\mathrm{O}(T(M))$. (We use the metric tensor to raise and lower indices and identify $T(M)=$ $T^{*}(M)$ ). Since these relations define an $\mathrm{O}(T(M))$ invariant subspace of $\bigotimes^{2 n}(T(M))$, we extend $P$ to be zero on the orthogonal complement of $\bigotimes^{n / 2}(V)$ in $\bigotimes^{2 n}(T(M))$ to extend $P$ to an equivariant action on the whole tensor algebra. Consequently, we can use H. Weyl's theorem to span $\mathcal{P}_{n, n, n}$ by expressions in which we alternate $n$ indices and contract in pairs the remaining $n$ indices.

For example, we compute:

$$
p_{1}=C \cdot R_{i j a b} R_{i j c d} d x^{a} \wedge d x^{b} \wedge d x^{c} \wedge d x^{d}
$$

for a suitable normalizing constant $C$ represents the first Pontrjagin form. In general, we will use letters $a, b, c, \ldots$ for indices to alternate on and indices $i, j, k, \ldots$ for indices to contract on. We let $P$ be such an element given by H. Weyl's theorem. There are some possibilities we can eliminate on a priori grounds. The Bianchi identity states:

$$
R_{i a b c} d x^{a} \wedge d x^{b} \wedge d x^{c}=\frac{1}{3}\left(R_{i a b c}+R_{i b c a}+R_{i c a b}\right) d x^{a} \wedge d x^{b} \wedge d x^{c}=0
$$

so that three indices of alternation never appear in any $R$... variable. Since there are $n / 2 R$ variables and $n$ indices of alternation, this implies each $R$ variable contains exactly two indices of alternation. We use the Bianchi identity again to express:

$$
\begin{aligned}
R_{i a j b} d x^{a} \wedge d x^{b} & =\frac{1}{2}\left(R_{i a j b}-R_{i b j a}\right) d x^{a} \wedge d x^{b} \\
& =\frac{1}{2}\left(R_{i a j b}+R_{i b a j}\right) d x^{a} \wedge d x^{b} \\
& =-\frac{1}{2} R_{i j b a} d x^{a} \wedge d x^{b}=\frac{1}{2} R_{i j a b} d x^{a} \wedge d x^{b} .
\end{aligned}
$$

This together with the other curvature identities means we can express $P$ in terms of $R_{i j a b} d x^{a} \wedge d x^{b}=\Omega_{i j}$ variables. Thus $P=P\left(\Omega_{i j}\right)$ is a polynomial in the components of the curvature matrix where we regard $\Omega_{i j} \in \Lambda^{2}\left(T^{*} M\right)$. (This differs by various factors of 2 from our previous definitions, but this is irrelevant to the present argument). $P(\Omega)$ is an $\mathrm{O}(n)$ invariant polynomial and thus is a real characteristic form. This completes the proof.

The remainder of this subsection is devoted to giving a combinatorial proof of this lemma in the case $n=m=p$ independent of H . Weyl's theorem. Since the Pontrjagin forms span a subspace of dimension $\pi(k)$ if $n=4 k$ we must show $\mathcal{P}_{n, n, n}=0$ if $n$ is not divisible by 4 and that $\operatorname{dim} \mathcal{P}_{4 k, 4 k, 4 k} \leq \pi(k)$ since then equality must hold.

We showed $n=2 j$ is even and any polynomial depends on the 2 -jets of the metric. We improve Lemma 2.5.1 as follows:

Lemma 2.5.8. Let $P$ satisfy the hypothesis of Lemma 2.5.1 and be a polynomial in the 2-jets of the metric. Then:
(a) Let $A=g_{12 / \alpha} A_{0}$ be a monomial of $P$. Either by interchanging 1 and 2 indices or by changing two 2 indices to 1 indices we can construct a monomial $A_{1}$ of the form $A_{1}=g_{11 / \beta} A_{0}^{\prime}$ which is a monomial of $P$.
(b) Let $A=g_{i j / 12} A_{0}$ be a monomial of $P$. Either by interchanging 1 and 2 indices or by changing two 2 indices to 1 indices we can construct a monomial of $A_{1}$ of the form $A_{1}=g_{i^{\prime} j^{\prime} / 11} A_{0}^{\prime}$ which is a monomial of $P$.
(c) The monomial $A_{1} \neq A$. If $\operatorname{deg}_{1}\left(A_{1}\right)=\operatorname{deg}_{1}(A)+2$ then $c(A, P) c\left(A_{1}, P\right)$ $>0$. Otherwise $\operatorname{deg}_{1}\left(A_{1}\right) \neq \operatorname{deg}_{1}(A)$ and $c(A, P) c\left(A_{1}, P\right)<0$.

Proof: We shall prove (a) as (b) is the same; we will also verify (c). Let $B=g_{11 / \alpha} A_{0}$ so $c(B, A(2 \rightarrow 1)) \neq 0$. Then $\operatorname{deg}_{1}(B)=\operatorname{deg}_{1}(A)+1$. Apply Lemma 2.5.2 to find $A_{1} \neq A$ so $c\left(B, A_{1}^{(1)}\right) \neq 0$. We noted earlier in the proof of 2.5 .1 that $A_{1}$ must have the desired form. If $c\left(A_{1}, B(2 \rightarrow 1)\right) \neq 0$ then $\operatorname{deg}_{1}\left(A_{1}\right)=\operatorname{deg}_{1}(B)+1=\operatorname{deg}_{1}(A)+2$ so $A_{1}$ is constructed from $A$ by changing two 2 to 1 indices. Furthermore, $c\left(B, A_{1}^{(1)}\right)>0$ and $c\left(B, A^{(1)}\right)<$ 0 implies $c(A, P) c\left(A_{1}, P\right)>0$. If, on the other hand, $c\left(A_{1}, B(1 \rightarrow 2)\right) \neq 0$ then $\operatorname{deg}_{1}\left(A_{1}\right)=\operatorname{deg}_{1}(B)-1=\operatorname{deg}_{1}(A)$ and $A$ changes to $A_{1}$ by interchanging a 2 and a 1 index. Furthermore, $c\left(B, A_{1}^{(1)}\right)<0$ and $c\left(B, A^{(1)}\right)<0$ implies $c(A, P) c\left(A_{1}, P\right)<0$ which completes the proof.

Let $P \in \mathcal{P}_{n, n, n}$ and express $P=P^{\prime} d x^{1} \wedge \cdots \wedge d x^{n} . \quad P^{\prime}=* P$ is a scalar invariant which changes sign if the orientation of the local coordinate system is reversed. We identify $P$ with $P^{\prime}$ for notational convenience and henceforth regard $P$ as a skew-invariant scalar polynomial. Thus $\operatorname{deg}_{k}(A)$ is odd for every $k$ and every monomial $A$ of $P$.

The indices with $\operatorname{deg}_{k}(A)=1$ play a particularly important role in our discussion. We say that an index $i$ touches an index $j$ in the monomial $A$
if $A$ is divisible by a variable $g_{i j / . .}$ or $g_{. . / i j}$ where ".." indicate two indices which are not of interest.

Lemma 2.5.9. Let $P \in \mathcal{P}_{n, n, n}$ with $P \neq 0$. Then there exists a monomial $A$ of $P$ so
(a) $\operatorname{deg}_{k}(A) \leq 3$ all $k$.
(b) If $\operatorname{deg}_{k}(A)=1$ there exists an index $j(k)$ which touches $k$ in $A$. The index $j(k)$ also touches itself in $A$.
(c) Let $\operatorname{deg}_{j} A+\operatorname{deg}_{k} A=4$. Suppose the index $j$ touches itself and the index $k$ in $A$. There is a unique monomial $A_{1}$ different from $A$ which can be formed from $A$ by interchanging a $j$ and $k$ index. $\operatorname{deg}_{j} A_{1}+\operatorname{deg}_{k} A_{1}=4$ and the index $j$ touches itself and the index $k$ in $A_{1} . c(A, P)+c\left(A_{1}, P\right)=0$.

Proof: Choose $A$ so the number of indices with $\operatorname{deg}_{k}(A)=1$ is minimal. Among all such choices, we choose $A$ so the number of indices which touch themselves is maximal. Let $\operatorname{deg}_{k}(A)=1$ then $k$ touches some index $j=$ $j(k) \neq k$ in $A$. Suppose $A$ has the form $A=g_{j k / . . ~}^{A_{0}}$ as the other case is similar. Suppose first that $\operatorname{deg}_{j}(A) \geq 5$. Use lemma 2.5.8 to find $A_{1}=$ $g_{k k / . .} A_{0}$. Then $\operatorname{deg}_{k}\left(A_{1}\right) \neq \operatorname{deg}_{k}(A)$ implies $\operatorname{deg}_{k}\left(A_{1}\right)=\operatorname{deg}_{k}(A)+2=3$. Also $\operatorname{deg}_{j}\left(A_{1}\right)=\operatorname{deg}_{j}(A)-2 \geq 3$ so $A_{1}$ has one less index with degree 1 which contradicts the minimality of the choice of $A$. We suppose next $\operatorname{deg}_{j}(A)=1$. If $T$ is the coordinate transformation interchanging $x_{j}$ and $x_{k}$, then $T$ reverses the orientation so $T^{*} P=-P$. However $\operatorname{deg}_{j}(A)=$ $\operatorname{deg}_{k}(A)=1$ implies $T^{*} A=A$ which contradicts the assumption that $A$ is a monomial of $P$. Thus $\operatorname{deg}_{j}(A)=3$ which proves (a).

Suppose $j$ does not touch itself in $A$. We use Lemma 2.5.8 to construct $A_{1}=g_{k k / . .} A_{0}^{\prime}$. Then $\operatorname{deg}_{k}\left(A_{1}\right)=\operatorname{deg}_{k}(A)+2=3$ and $\operatorname{deg}_{j}\left(A_{1}\right)=$ $\operatorname{deg}_{j}(A)-2=1$. This is a monomial with the same number of indices of degree 1 but which has one more index (namely $k$ ) which touches itself. This contradicts the maximality of $A$ and completes the proof of (b).

Finally, let $\mathcal{A}_{\nu}=\left\{k: \operatorname{deg}_{k}(A)=\nu\right\}$ for $\nu=1,3$. The map $k \mapsto j(k)$ defines an injective map from $\mathcal{A}_{1} \rightarrow \mathcal{A}_{3}$ since no index of degree 3 can touch two indices of degree 1 as well as touching itself. The equalities:
$n=\operatorname{card}\left(\mathcal{A}_{1}\right)+\operatorname{card}\left(\mathcal{A}_{3}\right) \quad$ and $\quad 2 n=\sum_{k} \operatorname{deg}_{k}(A)=\operatorname{card}\left(\mathcal{A}_{1}\right)+3 \operatorname{card}\left(\mathcal{A}_{3}\right)$
imply $2 \operatorname{card}\left(\mathcal{A}_{3}\right)=n$. Thus $\operatorname{card}\left(\mathcal{A}_{1}\right)=\operatorname{card}\left(\mathcal{A}_{3}\right)=n / 2$ and the map $k \mapsto j(k)$ is bijective in this situation.
(c) follows from Lemma 2.5.8 where $j=1$ and $k=2$. Since $\operatorname{deg}_{k}(A)=1$, $A_{1}$ cannot be formed by transforming two $k$ indices to $j$ indices so $A_{1}$ must be the unique monomial different from $A$ obtained by interchanging these indices. For example, if $A=g_{j j / a b} g_{j k / c d} A_{0}$, then $A_{1}=g_{j k / a b} g_{j j / c d} A_{0}$. The multiplicities involved are all 1 so we can conclude $c(A, P)+c\left(A_{1}, P\right)=$ 0 and not just $c(A, P) c\left(A_{1}, P\right)<0$.

Before further normalizing the choice of the monomial $A$, we must prove a lemma which relies on cubic changes of coordinates:

Lemma 2.5.10. Let $P$ be a polynomial in the 2-jets of the metric invariant under changes of coordinates of the form $y_{i}=x_{i}+c_{j k l} x_{j} x_{k} x_{l}$. Then:
(a) $g_{12 / 11}$ and $g_{11 / 12}$ divide no monomial $A$ of $P$.
(b) Let $A=g_{23 / 11} A_{0}$ and $B=g_{11 / 23} A_{0}$ then $A$ is a monomial of $P$ if and only if $B$ is a monomial of $P$. Furthermore, $c(A, P)$ and $c(B, P)$ have the same sign.
(c) Let $A=g_{11 / 22} A_{0}$ and $B=g_{22 / 11} A_{0}$, then $A$ is a monomial of $P$ if and only if $B$ is a monomial of $P$. Furthermore, $c(A, P)$ and $c(B, P)$ have the same sign.

Proof: We remark the use of the indices $1,2,3$ is for notational convenience only and this lemma holds true for any triple of distinct indices. Since $g_{i j}(X, G)=\delta_{i j}+O\left(x^{2}\right), g^{i j} / k l(X, G)\left(x_{0}\right)=-g_{i j / k l}(X, G)\left(x_{0}\right)$. Furthermore, under changes of this sort, $d_{y}^{\alpha}\left(x_{0}\right)=d_{x}^{\alpha}\left(x_{0}\right)$ if $|\alpha|=1$. We consider the change of coordinates:

$$
\begin{array}{rlrlrl}
y_{2} & =x_{2}+c x_{1}^{3}, & & y_{k} & =k_{x} & \\
\text { otherwise } \\
d y_{2} & =d x_{2}+3 c x_{1}^{2} d x_{1}, & d y_{k} & =d x_{k} & & \text { otherwise }
\end{array}
$$

with

$$
\begin{aligned}
g^{12}(Y, G) & =g^{12}(X, G)+3 c x_{1}^{2}+O\left(x^{4}\right), \\
g^{i j}(Y, G) & =g^{i j}(X, G)+O\left(x^{4}\right) \text { otherwise } \\
g_{12 / 11}(Y, G)\left(x_{0}\right) & =g_{12 / 11}(X, G)\left(x_{0}\right)-6 c \\
g_{i j / k l}(Y, G)\left(x_{0}\right) & =g_{i j / k l}(X, G)\left(x_{0}\right) \text { otherwise. }
\end{aligned}
$$

We decompose $A=\left(g_{12 / 11}\right)^{\nu} A_{0}$. If $\nu>0, T^{*}(A)=A-6 \nu c\left(g_{12 / 11}\right)^{\nu-1} A_{0}+$ $O\left(c^{2}\right)$. Since $T^{*}(P)=P$, and since there is no way to cancel this additional contribution, $A$ cannot be a monomial of $P$ so $g_{12 / 11}$ divides no monomial of $P$.

Next we consider the change of coordinates:

$$
\begin{aligned}
y_{1} & =x_{1}+c x_{1}^{2} x_{2}, & y_{k} & =x_{k} & & \text { otherwise } \\
d y_{1} & =d x_{1}+2 c x_{1} x_{2} d x_{1}+c x_{1}^{2} d x_{2}, & d y_{k} & =d x_{k} & & \text { otherwise }
\end{aligned}
$$

with

$$
\begin{array}{rlrl}
g_{11 / 12}(Y, G)\left(x_{0}\right) & =g_{11 / 12}(X, G)\left(x_{0}\right)-4 c \\
g_{12 / 11}(Y, G)\left(x_{0}\right) & =g_{12 / 11}(X, G)\left(x_{0}\right)-2 c \\
g_{i j / k l}(Y, G)\left(x_{0}\right) & =g_{i j / k l}(X, G)\left(x_{0}\right) & \text { otherwise. }
\end{array}
$$

We noted $g_{12 / 11}$ divides no monomial of $P$. If $A=\left(g_{11 / 12}\right)^{\nu} A_{0}$, then $\nu>0$ implies $T^{*}(A)=A-4 c v\left(g_{11 / 12}\right)^{\nu-1} A_{0}+O\left(c^{2}\right)$. Since there would be no
way to cancel such a contribution, $A$ cannot be a monomial of $P$ which completes the proof of (a).

Let $A_{0}^{\prime}$ be a monomial not divisible by any of the variables $g_{23 / 11}, g_{12 / 13}$, $g_{13 / 12}, g_{11 / 23}$ and let $A(p, q, r, s)=\left(g_{23 / 11}\right)^{p}\left(g_{12 / 13}\right)^{q}\left(g_{13 / 12}\right)^{r}\left(g_{11 / 23}\right)^{s} A_{0}^{\prime}$. We set $c(p, q, r, s)=c(A(p, q, r, s), P)$ and prove (b) by establishing some relations among these coefficients. We first consider the change of coordinates,

$$
\begin{aligned}
y_{2} & =x_{2}+c x_{1}^{2} x_{3}, & y_{k} & =x_{k} & & \text { otherwise } \\
d y_{2} & =d x_{2}+2 c x_{1} x_{3} d x_{1}+c x_{1}^{2} d x_{3}, & d y_{k} & =d x_{k} & & \text { otherwise }
\end{aligned}
$$

with

$$
\begin{aligned}
g_{12 / 13}(Y, G)\left(x_{0}\right) & =g_{12 / 13}(X, G)\left(x_{0}\right)-2 c \\
g_{23 / 11}(Y, G)\left(x_{0}\right) & =g_{23 / 11}(X, G)\left(x_{0}\right)-2 c \\
g_{i j / k l}(Y, G)\left(x_{0}\right) & =g_{i j / k l}(X, G)\left(x_{0}\right) \quad \text { otherwise. }
\end{aligned}
$$

We compute that:

$$
\begin{aligned}
& T^{*}(A(p, q, r, s)) \\
& \quad=A(p, q, r, s)+c\{-2 q A(p, q-1, r, s)-2 p A(p-1, q, r, s)\}+O\left(c^{2}\right)
\end{aligned}
$$

Since $T^{*}(P)=P$ is invariant, we conclude

$$
p c(p, q, r, s)+(q+1) c(p-1, q+1, r, s)=0
$$

By interchanging the roles of 2 and 3 in the argument we also conclude:

$$
p c(p, q, r, s)+(r+1) c(p-1, q, r+1, s)=0
$$

(We set $c(p, q, r, s)=0$ if any of these integers is negative.)
Next we consider the change of coordinates:

$$
\begin{aligned}
y_{1} & =x_{1}+c x_{1} x_{2} x_{3}, & & y_{k}
\end{aligned}=x_{k} \quad \begin{array}{ll}
\text { otherwise } \\
d y_{1} & =d x_{1}+c x_{1} x_{2} d x_{3}+c x_{1} x_{3} d x_{2}+c x_{2} x_{3} d x_{1}, \\
d y_{k} & =d x_{k}
\end{array} \begin{array}{ll}
\text { otherwise }
\end{array}
$$

with

$$
\begin{aligned}
& g_{11 / 23}(Y, G)\left(x_{0}\right)=g_{11 / 23}(X, G)\left(x_{0}\right)-2 c \\
& g_{12 / 13}(Y, G)\left(x_{0}\right)=g_{12 / 13}(X, G)\left(x_{0}\right)-c \\
& g_{13 / 12}(Y, G)\left(x_{0}\right)=g_{13 / 12}(Y, G)\left(x_{0}\right)-c
\end{aligned}
$$

so that

$$
\begin{aligned}
T^{*}(A(p, q, r, s))= & A(p, q, r, s)+c\{-2 s A(p, q, r, s-1) \\
& -r A(p, q, r-1, s)-q A(p, q-1, r, s)\}+O\left(c^{2}\right)
\end{aligned}
$$

This yields the identities

$$
(q+1) c(p, q+1, r, s-1)+(r+1) c(p, q, r+1, s-1)+2 s c(p, q, r, s)=0
$$

Let $p \neq 0$ and let $A=g_{23 / 11} A_{0}^{\prime}$ be a monomial of $P$. Then $c(p-1, q+1, r, s)$ and $c(p-1, q, r+1, s)$ are non-zero and have the opposite sign as $c(p, q, r, s)$. Therefore $(q+1) c(p-1, q+1, r, s)+(r+1) c(p-1, q, r+1, s)$ is non-zero which implies $c(p-1, q, r, s+1)(s+1)$ is non-zero and has the same sign as $c(p, q, r, s)$. This shows $g_{11 / 23} A_{0}^{\prime}$ is a monomial of $P$. Conversely, if $A$ is not a monomial of $P$, the same argument shows $g_{11 / 23} A_{0}^{\prime}$ is not a monomial of $P$. This completes the proof of (b).

The proof of (c) is essentially the same. Let $A_{0}^{\prime}$ be a monomial not divisible by the variables $\left\{g_{11 / 22}, g_{12 / 12}, g_{22 / 11}\right\}$ and let

$$
A(p, q, r)=\left(g_{11 / 22}\right)^{p}\left(g_{12 / 12}\right)^{q}\left(g_{22 / 11}\right)^{r}
$$

Let $c(p, q, r)=c(A(p, q, r), P)$. Consider the change of coordinates:

$$
\begin{aligned}
y_{1} & =x_{1}+c x_{1} x_{2}^{2}, & y_{k} & =x_{k} & & \text { otherwise } \\
d y_{1} & =d x_{1}+c x_{2}^{2} d x_{1}+2 c x_{1} x_{2} d x_{2}, & d y_{k} & =d x_{k} & & \text { otherwise }
\end{aligned}
$$

with

$$
\begin{array}{rlrl}
g_{11 / 22}(Y, g)\left(x_{0}\right) & =g_{11 / 22}(X, G)\left(x_{0}\right)-4 c \\
g_{12 / 12}(Y, G)\left(x_{0}\right) & =g_{12 / 12}(X, G)\left(x_{0}\right)-2 c \\
g_{i j / k l}(Y, G)\left(x_{0}\right) & =g_{i j / k l}(X, G)\left(x_{0}\right) & \text { otherwise. }
\end{array}
$$

This yields the relation $2 p c(p, q, r)+(q+1) c(p-1, q+1, r)=0$. By interchanging the roles of 1 and 2 we obtain the relation $2 r c(p, q, r)+(q+$ 1) $c(p, q+1, r-1)=0$ from which (c) follows.

This step in the argument is functionally equivalent to the use made of the $\left\{R_{i j k l}\right\}$ variables in the argument given previously which used H . Weyl's formula. It makes use in an essential way of the invariance of $P$ under a wider group than just first and second order transformations. For the Euler form, by contrast, we only needed first and second order coordinate transformations.

We can now construct classifying monomials using these lemmas. Fix $n=2 n_{1}$ and let $A$ be the monomial of $P$ given by Lemma 2.5.9. By making a coordinate permutation, we may assume $\operatorname{deg}_{k}(A)=3$ for $k \leq n_{1}$ and $\operatorname{deg}_{k}(A)=1$ for $k>n_{1}$. Let $x(i)=i+n_{1}$ for $1 \leq i \leq n_{0}$; we may assume the index $I$ touches itself and $x(i)$ in $A$ for $i \leq n_{0}$.

We further normalize the choice of $A$ as follows. Either $g_{11 / i j}$ or $g_{i j / 11}$ divides $A$. Since $\operatorname{deg}_{1}(A)=3$ this term is not $g_{11 / 11}$ and by Lemma 2.5.10 it is not $g_{11 / 1 x(1)}$ nor $g_{1 x(1) / 11}$. By Lemma 2.5.10(b) or 2.5.10(c), we may assume $A=g_{11 / i j} \ldots$ for $i, j \geq 2$. Since not both $i$ and $j$ can have degree

1 in $A$, by making a coordinate permutation if necessary we may assume that $i=2$. If $j=2$, we apply Lemma 2.5.9(c) to the indices 2 and $x(2)$ to perform an interchange and assume $A=g_{11 / 2 x(2)} A_{0}$. The index 2 must touch itself elsewhere in $A$. We apply the same considerations to choose $A$ in the form $A=g_{11 / 2 x(2)} g_{22 / i j} A_{0}$. If $i$ or $j$ is 1 , the cycle closes and we express $A=g_{11 / 2 x(2)} g_{11 / 1 x(1)} A_{0}$ where $\operatorname{deg}_{k}\left(A_{0}\right)=0$ for $k=1,2,1+n_{1}$, $2+n_{1}$. If $i, j>2$ we continue this argument until the cycle closes. This permits us to choose $A$ to be a monomial of $P$ in the form:

$$
A=g_{11 / 2 x(2)} g_{22 / 3 x(3)} \ldots g_{j-1, j-1 / j x(j)} g_{j j / 1 x(1)} A_{0}
$$

where $\operatorname{deg}_{k}\left(A_{0}\right)=0$ for $1 \leq k \leq j$ and $n_{1}+1 \leq k \leq n_{1}+j$.
We wish to show that the length of the cycle involved is even. We apply Lemma 2.5.10 to show

$$
B=g_{2 x(2) / 11} g_{3 x(3) / 22} \ldots g_{1 x(1) / j j} A_{0}
$$

satisfies $c(A, P) c(B, P)>0$. We apply Lemma 2.5.9 a total of $j$ times to see

$$
C=g_{22,1 x(1)} g_{33 / 2 x(2)} \ldots g_{11 / j x(j)} A_{0}
$$

satisfies $c(B, P) c(C, P)(-1)^{j}>0$. We now consider the even permutation:

$$
\begin{aligned}
& \rho(1)=j \\
& \rho(k)= \begin{cases}k-1, & 2 \leq k \leq j \\
k, & j<k \leq n_{1}\end{cases} \\
& \rho(k(j))=x(\rho(j)) \quad \text { for } 1 \leq j \leq n_{1}
\end{aligned}
$$

to see that $c\left(C_{\rho}, P\right) c(C, P)>0$. However $A=C_{\rho}$ so $c(A, P)^{2}(-1)^{j}>0$ which shows $j$ is necessarily even. (This step is formally equivalent to using the skew symmetry of the curvature tensor to show that only polynomials of even degree can appear to give non-zero real characteristic forms).

We decompose $A_{0}$ into cycles to construct $A$ inductively so that $A$ has the form:

$$
\begin{aligned}
A=\{ & \left.g_{11 / 2 x(2)} \ldots g_{i_{1} i_{1} / 1 x(1)}\right\} \\
& \left\{g_{i_{1}+1, i_{1}+1 / i_{1}+2, x\left(i_{1}+2\right)} \ldots g_{i_{1}+i_{2}, i_{1}+i_{2} / i_{1}, x\left(i_{1}\right)}\right\} \ldots
\end{aligned}
$$

where we decompose $A$ into cycles of length $i_{1}, i_{2}, \ldots$ with $n_{1}=i_{1}+\cdots+i_{j}$. Since all the cycles must have even length, $\ell(A)=n / 2$ is even so $n$ is divisible by 4 .

We let $n=4 k$ and let $\rho$ be a partition of $k=k_{1}+\cdots+k_{j}$. We let $A_{\rho}$ be defined using the above equation where $i_{1}=2 k_{1}, i_{2}=2 k_{2}, \ldots$. By making a coordinate permutation we can assume $i_{1} \geq i_{2} \geq \cdots$. We have shown
that if $P=0$, then $c\left(A_{\rho}, P\right)=0$ for some $\rho$. Since there are exactly $\pi(k)$ such partitions, we have constructed a family of $\pi(k)$ linear functionals on $\mathcal{P}_{n, n, n}$ which form a seperating family. This implies $\operatorname{dim} \mathcal{P}_{n, n, n}<\pi(k)$ which completes the proof.

We conclude this subsection with a few remarks on the proofs we have given of Theorem 2.4.7 and Lemma 2.5.6. We know of no other proof of Theorem 2.4.7 other than the one we have given. H. Weyl's theorem is only used to prove the surjectivity of the restriction map $r$ and is inessential in the axiomatic characterization of the Euler form. This theorem gives an immediate proof of the Gauss-Bonnet theorem using heat equation methods. It is also an essential step in settling Singer's conjecture for the Euler form as we shall discuss in the fourth chapter. The fact that $r\left(E_{m}\right)=0$ is, of course, just an invariant statement of the fact $E_{m}$ is an unstable characteristic class; this makes it difficult to get hold of axiomatically in contrast to the Pontrjagin forms which are stable characteristic classes.

We have discussed both of the proofs of Lemma 2.5.6 which exist in the literature. The proof based on H. Weyl's theorem and on geodesic normal coordinates is more elegant, but relies heavily on fairly sophisticated theorems. The second is the original proof and is more combinatorial. It is entirely self-contained and this is the proof which generalizes to Kaehler geometry to yield an axiomatic characterization of the Chern forms of $T_{c}(M)$ for a holomorphic Kaehler manifold. We shall discuss this case in more detail in section 3.7.

### 2.6. Invariance Theory and Mixed Characteristic Classes of the Tangent Space and of a Coefficient Bundle.

In the previous subsection, we gave in Lemma 2.5.6 an axiomatic characterization of the Pontrjagin forms in terms of functorial properties. In discussing the Hirzebruch signature formula in the next chapter, it will be convenient to have a generalization of this result to include invariants which also depend on the derivatives of the connection form on an auxilary bundle.

Let $V$ be a complex vector bundle. We assume $V$ is equipped with a Hermitian fiber metric and let $\nabla$ be a Riemannian or unitary connection on $V$. We let $\vec{s}=\left(s_{1}, \ldots, s_{a}, \ldots, s_{v}\right)$ be a local orthonormal frame for $V$ and introduce variables $\omega_{a b i}$ for the connection 1-form;

$$
\nabla\left(s_{a}\right)=\omega_{a b i} d x^{i} \otimes s_{b}, \quad \text { i.e., } \nabla \vec{s}=\omega \otimes \vec{s} .
$$

We introduce variables $\omega_{a b i / \alpha}=d_{x}^{\alpha}\left(\omega_{a b i}\right)$ for the partial derivatives of the connection 1-form. We shall also use the notation $\omega_{a b i / j k \ldots}$. We use indices $1 \leq a, b, \cdots \leq v$ to index the frame for $V$ and indices $1 \leq i, j, k \leq m$ for the tangent space variables. We define:

$$
\operatorname{ord}\left(\omega_{a b i / \alpha}\right)=1+|\alpha| \quad \text { and } \quad \operatorname{deg}_{k}\left(\omega_{a b i / \alpha}\right)=\delta_{i, k}+\alpha(k)
$$

We let $\mathcal{Q}$ be the polynomial algebra in the $\left\{\omega_{a b i / \alpha}\right\}$ variables for $|\alpha| \geq 1$; if $Q \in \mathcal{Q}$ we define the evaluation $Q(X, \vec{s}, \nabla)\left(x_{0}\right)$. We normalize the choice of frame $\vec{s}$ by requiring $\nabla(\vec{s})\left(x_{0}\right)=0$. We also normalize the coordinate system $X$ as before so $X\left(x_{0}\right)=0, g_{i j}(X, G)\left(x_{0}\right)=\delta_{i j}$, and $g_{i j / k}(X, G)\left(x_{0}\right)=$ 0 . We say $Q$ is invariant if $Q(X, \vec{s}, \nabla)\left(x_{0}\right)=Q\left(Y, \vec{s}^{\prime}, \nabla\right)\left(x_{0}\right)$ for any normalized frames $\vec{s}, \vec{s}^{\prime}$ and normalized coordinate systems $X, Y$; we denote this common value by $Q(\nabla)$ (although it also depends in principle on the metric tensor and the 1-jets of the metric tensor through our normalization of the coordinate system $X$ ). We let $\mathcal{Q}_{m, p, v}$ denote the space of all invariant $p$-form valued polynomials in the $\left\{\omega_{a b i / \alpha}\right\}$ variables for $|\alpha| \geq 1$ defined on a manifold of dimension $m$ and for a vector bundle of complex fiber dimension $v$. We let $\mathcal{Q}_{m, n, p, v}$ denote the subspace of invariant polynomials homogeneous of order $n$ in the jets of the connection form. Exactly as was done for the $\mathcal{P}_{*}$ algebra in the jets of the metric, we can show there is a direct sum decomposition

$$
\mathcal{Q}_{m, p, v}=\bigoplus_{n} \mathcal{Q}_{m, n, p, v} \quad \text { and } \quad \mathcal{Q}_{m, n, p, v}=0 \quad \text { for } n-p \text { odd }
$$

Let $Q\left(g A g^{-1}\right)=Q(A)$ be an invariant polynomial of order $q$ in the components of a $v \times v$ matrix. Then $Q(\Omega)$ defines an element of $\mathcal{Q}_{m, 2 q, q, v}$ for
any $m \geq 2 q$. By taking $Q \cdot \bar{Q}$ we can define scalar valued invariants and by taking $\delta(Q)$ we can define other form valued invariants in $\mathcal{Q}_{m, 2 q+1,2 q-1, v}$. Thus there are a great many such form valued invariants.

In addition to this algebra, we let $\mathcal{R}_{m, n, p, v}$ denote the space of $p$-form valued invariants which are homogeneous of order $n$ in the $\left\{g_{i j / \alpha}, \omega_{a b k / \beta}\right\}$ variables for $|\alpha| \geq 2$ and $|\beta| \geq 1$. The spaces $\mathcal{P}_{m, n, p}$ and $\mathcal{Q}_{m, n, p, v}$ are both subspaces of $\mathcal{R}_{m, n, p, v}$. Furthermore wedge product gives a natural map $\mathcal{P}_{m, n, p} \otimes \mathcal{Q}_{m, n^{\prime}, p^{\prime}, v} \rightarrow \mathcal{R}_{m, n+n^{\prime}, p+p^{\prime}, v}$. We say that $R \in \mathcal{R}_{m, p, p, v}$ is a characteristic form if it is in the linear span of wedge products of Pontrjagin forms of $T(M)$ and Chern forms of $V$. The characteristic forms are characterized abstractly by the following Theorem. This is the generalization of Lemma 2.5.6 which we shall need in discussing the signature and spin complexes.

Theorem 2.6.1.
(a) $\mathcal{R}_{m, n, p, v}=0$ if $n<p$ or if $n=p$ and $n$ is odd.
(b) If $R \in \mathcal{R}_{m, n, n, v}$ then $R$ is a characteristic form.

Proof: The proof of this fact relies heavily on Lemma 2.5.6 but is much easier. We first need the following generalization of Lemma 2.5.1:

Lemma 2.6.2. Using the notation of Lemma 2.5 .1 we define $T_{z}^{*}(R)$ if $R$ is a scalar invariant in the $\left\{g_{i j / \alpha}, \omega_{a b j / \beta}\right\}$ variables.
(a) Let $g_{12 / \alpha}$ divide some monomial of $R$, then $g_{\text {../ } / \beta}$ divides some other monomial of $R$.
(b) Let $g_{. . / \alpha}$ divide some monomial of $R$, then $g_{\text {../ } / \beta}$ divides some other monomial of $R$ where $\beta(1)=\alpha(1)+\alpha(2)$ and $\beta(2)=0$.
(c) Let $\omega_{a b i / \alpha}$ divide some monomial of $R$ for $i>2$, then $\omega_{a b i / \beta}$ divides some other monomial of $R$ where $\beta(1)=\alpha(1)+\alpha(2)$ and $\beta(2)=0$.

The proof of this is exactly the same as that given for Lemma 2.5.1 and is therefore omitted.

The proof of Theorem 2.6.1 parallels the proof of Lemma 2.5.6 for a while so we summarize the argument briefly. Let $0 \neq R \in \mathcal{R}_{m, n, p, v}$. The same argument given in Lemma 2.5.3 shows an invariant polynomial is homogeneous of order $n$ if $R\left(c^{2} G, \nabla\right)=c^{p-n} R(G, \nabla)$ which gives a invariant definition of the order of a polynomial. The same argument as given in Lemma 2.5.4 shows $n-p$ must be even and that if $A$ is a monomial of $R, \quad A$ is a monomial of exactly one of the $R_{I}$. If $c\left(A, R_{I}\right)=0$ then $\operatorname{deg}_{k}(A)+\operatorname{deg}_{k}(I)$ is always even. We decompose $A$ in the form:

$$
A=g_{i_{1} j_{1} / \alpha_{1}} \ldots g_{i_{q} j_{q} / \alpha_{q}} \omega_{a_{1} b_{1} k_{1} / \beta_{1}} \ldots \omega_{a_{r} b_{r} k_{r} / \beta_{r}}=A^{g} A^{\omega}
$$

and define $\ell(A)=q+r$ to the length of $A$. We argue using Lemma 2.6.2 to choose $A$ so $\operatorname{deg}_{k}\left(A^{g}\right)=0$ for $k>2 q$. By making a coordinate permutation we can assume that the $k_{\nu} \leq 2 q+r$ for $1 \leq \nu \leq r$. We
apply Lemma 2.6.2(c) a total of $r$ times to choose the $\beta_{i}$ so $\beta_{1}(k)=0$ for $k>2 q+r+1, \beta_{2}(k)=0$ for $k>2 q+r+2, \ldots, \beta_{r}(k)=0$ for $k>2 q+2 r$. This chooses $A$ so $\operatorname{deg}_{k}(A)=0$ for $k>2 \ell(A)$. If $A$ is a monomial of $R_{I}$ for $I=\left\{1 \leq i_{1}<\cdots<i_{p} \leq m\right\}$ then $\operatorname{deg}_{i_{p}}(A)$ is odd. We estimate $p \leq i_{p} \leq 2 \ell(A) \leq \sum\left|\alpha_{\nu}\right|+\sum\left(\left|\beta_{\mu}\right|+1\right)=n$ so that $\mathcal{P}_{m, n, p, v}=\{0\}$ if $n<p$ or if $n-p$ is odd which proves (a) of Theorem 2.6.1.

In the limiting case, we must have equalities so $\left|\alpha_{\nu}\right|=2$ and $\left|\beta_{\mu}\right|=1$. Furthermore, $i_{p}=p$ so there is some monomial $A$ so $\operatorname{deg}_{k}(A)=0$ for $k>p=n$ and $A d x^{1} \wedge \cdots \wedge d x^{p}$ appears in $R$. There is a natural restriction map

$$
r: \mathcal{R}_{m, n, p, v} \rightarrow \mathcal{R}_{m-1, n, p, v}
$$

and our argument shows $r: \mathcal{R}_{m, n, n, v} \rightarrow \mathcal{R}_{m-1, n, n, v}$ is injective for $n<m$. Since the restriction of a characteristic form is a characteristic form, it suffices to prove (b) of Theorem 2.6.1 for the case $m=n=p$.

Let $0 \neq R \in \mathcal{R}_{n, n, n, v}$ then $R$ is a polynomial in the $\left\{g_{i j / k l}, \omega_{a b i / j}\right\}$ variables. The restriction map $r$ was defined by considering products $M_{1} \times$ $S^{1}$ but there are other functorial constructions which give rise to useful projections. Fix non-negative integers $(s, t)$ so that $n=s+t$. Let $M_{1}$ be a Riemannian manifold of dimension $s$. Let $M_{2}$ be the flat torus of dimension $t$ and let $V_{2}$ be a vector bundle with connection $\nabla_{2}$ over $M_{2}$. Let $M=M_{1} \times M_{2}$ with the product metric and let $V$ be the natural extension of $V_{2}$ to $M$ which is flat in the $M_{1}$ variables. More exactly, if $\pi_{2}: M \rightarrow M_{2}$ is a projection on the second factor, then $(V, \nabla)=\pi_{2}^{*}\left(V_{2}, \nabla_{2}\right)$ is the pull back bundle with the pull back connection. We define

$$
\pi_{(s, t)}(R)\left(G_{1}, \nabla_{2}\right)=R\left(G_{1} \times 1, \nabla\right)
$$

Using the fact that $\mathcal{P}_{s, n_{1}, p_{1}}=0$ for $s<p_{1}$ or $n_{1}<p_{1}$ and the fact $\mathcal{Q}_{t, n_{2}, p_{2}, k}=0$ for $t<p_{2}$ or $n_{2}<p_{2}$ it follows that $\pi_{(s, t)}$ defines a map

$$
\pi_{(s, t)}: \mathcal{R}_{n, n, v} \rightarrow \mathcal{P}_{s, s, s} \otimes \mathcal{Q}_{t, t, t, v}
$$

More algebraically, let $A=A^{g} A^{\omega}$ be a monomial, then we define:

$$
\pi_{(s, t)}(A)= \begin{cases}0 & \text { if } \operatorname{deg}_{k}\left(A^{g}\right)>0 \text { for } k>s \text { or } \operatorname{deg}_{k}\left(A^{\omega}\right)>0 \text { for } k \leq s \\ A & \text { otherwise. }\end{cases}
$$

The only additional relations imposed are to set $g_{i j / k l}=0$ if any of these indices exceeds $s$ and to set $\omega_{a b i / j}=0$ if either $i$ or $j$ is less than or equal to $s$.

We use these projections to reduce the proof of Theorem 2.6.1 to the case in which $R \in \mathcal{Q}_{t, t, t, v}$. Let $0 \neq R \in \mathcal{R}_{n, n, n, v}$ and let $A=A^{g} A^{\omega}$ be a monomial of $R$. Let $s=2 \ell\left(A^{g}\right)=\operatorname{ord}\left(A^{g}\right)$ and let $t=n-s=2 \ell\left(A^{\omega}\right)=$
$\operatorname{ord}\left(A^{\omega}\right)$. We choose $A$ so $\operatorname{deg}_{k}\left(A^{g}\right)=0$ for $k>s$. Since $\operatorname{deg}_{k}(A) \geq 1$ must be odd for each index $k$, we can estimate:

$$
t \leq \sum_{k>s} \operatorname{deg}_{k}(A)=\sum_{k>s} \operatorname{deg}_{k}\left(A^{\omega}\right) \leq \sum_{k} \operatorname{deg}_{k}\left(A^{\omega}\right)=\operatorname{ord}\left(A^{\omega}\right)=t
$$

As all these inequalities must be equalities, we conclude $\operatorname{deg}_{k}\left(A^{\omega}\right)=0$ for $k \leq s$ and $\operatorname{deg}_{k}\left(A^{\omega}\right)=1$ for $k>s$. This shows in particular that $\pi_{(s, t)}(R) \neq 0$ for some $(s, t)$ so that

$$
\bigoplus_{s+t=n} \pi_{(s, t)}: \mathcal{R}_{n, n, n, v} \rightarrow \bigoplus_{s+t=n} \mathcal{P}_{s, s, s} \otimes \mathcal{Q}_{t, t, t, v}
$$

is injective.
We shall prove that $\mathcal{Q}_{t, t, t, v}$ consists of characteristic forms of $V$. We showed earlier that $\mathcal{P}_{s, s, s}$ consists of Pontrjagin forms of $T(M)$. The characteristic forms generated by the Pontrjagin forms of $T(M)$ and of $V$ are elements of $\mathcal{R}_{n, n, n, v}$ and $\pi_{(s, t)}$ just decomposes such products. Thererfore $\pi$ is surjective when restricted to the subspace of characteristic forms. This proves $\pi$ is bijective and also that $\mathcal{R}_{n, n, n, v}$ is the space of characteristic forms. This will complete the proof of Theorem 2.6.1.

We have reduced the proof of Theorem 2.6.1 to showing $\mathcal{Q}_{t, t, t, v}$ consists of the characteristic forms of $V$. We noted that $0 \neq Q \in \mathcal{Q}_{t, t, t, v}$ is a polynomial in the $\left\{\omega_{a b i / j}\right\}$ variables and that if $A$ is a monomial of $Q$, then $\operatorname{deg}_{k}(A)=1$ for $1 \leq k \leq t$. Since $\operatorname{ord}(A)=t$ is even, we conclude $\mathcal{Q}_{t, t, t, v}=0$ if $t$ is odd.

The components of the curvature tensor are given by:

$$
\Omega_{a b i j}=\omega_{a b i / j}-\omega_{a b i / j} \quad \text { and } \quad \Omega_{a b}=\sum \Omega_{a b i j} d x_{i} \wedge d x_{j}
$$

up to a possible sign convention and factor of $\frac{1}{2}$ which play no role in this discussion. If $A$ is a monomial of $P$, we decompose:

$$
A=\omega_{a_{1} b_{1} i_{1} / i_{2}} \ldots \omega_{a_{u} b_{u} i_{t-1} / i_{t}} \quad \text { where } 2 u=t
$$

All the indices $i_{\nu}$ are distinct. If $\rho$ is a permutation of these indices, then $c(A, P)=\operatorname{sign}(\rho) c\left(A^{\rho}, P\right)$. This implies we can express $P$ in terms of the expressions:

$$
\begin{gathered}
\bar{A}=\left(\omega_{a_{1} b_{1} i_{1} / i_{2}}-\omega_{a_{1} b_{1} i_{2} / i_{1}}\right) \ldots\left(\omega_{a_{u} b_{u} i_{t-1} / i_{t}}-\omega_{a_{u} b_{u} i_{t} / i_{t-1}}\right) \\
d x_{i_{1}} \wedge \cdots \wedge d x_{i_{t}} \\
=\Omega_{a_{1} b_{1} i_{1} i_{2}} \ldots \Omega_{a_{u} b_{u} i_{t-1} i_{t}} d x_{i_{1}} \wedge \cdots \wedge d x_{i_{t}}
\end{gathered}
$$

Again, using the alternating nature of these expression, we can express $P$ in terms of expressions of the form:

$$
\Omega_{a_{1} b_{1}} \wedge \cdots \wedge \Omega_{a_{u} b_{u}}
$$

so that $Q=Q(\Omega)$ is a polynomial in the components $\Omega_{a b}$ of the curvature. Since the value of $Q$ is independent of the frame chosen, $Q$ is the invariant under the action of $\mathrm{U}(v)$. Using the same argument as that given in the proof of Lemma 2.1.3 we see that in fact $Q$ is a characteristic form which completes the proof.

We conclude this subsection with some uniqueness theorems regarding the local formulas we have been considering. If $M$ is a holomorphic manifold of real dimension $m$ and if $V$ is a complex vector bundle of fiber dimension $v$, we let $\mathcal{R}_{m, m, p, v}^{c h}$ denote the space of $p$-form valued invariants generated by the Chern forms of $V$ and of $T_{c}(M)$. There is a suitable axiomatic characterization of these spaces using invariance theory for Kaehler manifolds which we shall discuss in section 3.7. The uniqueness result we shall need in proving the Hirzebruch signature theorem and the RiemannRoch theorem is the following:

Lemma 2.6.3.
(a) Let $0 \neq R \in \mathcal{R}_{m, m, m, v}$ then there exists $(M, V)$ so $M$ is oriented and $\int_{M} R(G, \nabla) \neq 0$.
(b) Let $0 \neq R \in \mathcal{R}_{m, m, m, v}^{c h}$ then there exists $(M, V)$ so $M$ is a holomorphic manifold and $\int_{M} R(G, \nabla) \neq 0$.
Proof: We prove (a) first. Let $\rho=\left\{1 \leq i_{1} \leq \cdots \leq i_{\nu}\right\}$ be a partition $k(\rho)=i_{1}+\cdots+i_{\nu}$ for $4 k=s \leq m$. Let $\left\{M_{\rho}^{r}\right\}$ be the collection of manifolds of dimension $s$ discussed in Lemma 2.3.4. Let $P_{\rho}$ be the corresponding real characteristic form so

$$
\int_{M_{\rho}^{r}} P_{\tau}(G)=\delta_{\rho, \tau}
$$

The $\left\{P_{\rho}\right\}$ forms a basis for $\mathcal{P}_{m, s, s}$ for any $s \leq m$. We decompose

$$
R=\sum_{\rho} P_{\rho} Q_{\rho} \quad \text { for } Q_{\rho} \in \mathcal{Q}_{m, t, t, k}, \text { where } t+s=m
$$

This decomposition is, of course, nothing but the decomposition defined by the projections $\pi_{(s, t)}$ discussed in the proof of the previous lemma.

Since $R \neq 0$, at least one of the $Q_{\rho} \neq 0$. We choose $\rho$ so $k(\rho)$ is maximal with $Q_{\rho} \neq 0$. We consider $M=M_{\rho}^{r} \times M_{2}$ and $(V, \nabla)=\pi_{2}^{*}\left(V_{2}, \nabla_{2}\right)$ where $\left(V_{2}, \nabla_{2}\right)$ is a bundle over $M_{2}$ which will be specified later. Then we compute:

$$
\int_{M} P_{\tau}(G) Q_{\tau}(\nabla)=0
$$

unless $\operatorname{ord}\left(Q_{\tau}\right) \leq \operatorname{dim}\left(M_{2}\right)$ since $\nabla$ is flat along $M_{\rho}^{r}$. This implies $k(\tau) \geq$ $k(\rho)$ so if this integral is non-zero $k(\tau)=k(\rho)$ by the maximality of $\rho$. This implies

$$
\int_{M} P_{\tau}(G) Q_{\tau}(\nabla)=\int_{M_{\rho}^{r}} P_{\tau}(G) \cdot \int_{M_{2}} Q_{\tau}\left(\nabla_{2}\right)=\delta_{\rho, \tau} \int_{M_{2}} Q_{\tau}\left(\nabla_{2}\right)
$$

This shows that

$$
\int_{M} R(G, \nabla)=\int_{M_{2}} Q_{\rho}\left(\nabla_{2}\right)
$$

and reduces the proof of this lemma to the special case $Q \in \mathcal{Q}_{m, m, m, v}$.
Let $A$ be a $v \times v$ complex matrix and let $\left\{x_{1}, \ldots, x_{v}\right\}$ be the normalized eigenvalues of $A$. If $2 k=m$ and if $\rho$ is a partition of $k$ we define

$$
x_{\rho}=x_{1}^{i_{1}} \ldots x_{\nu}^{i_{\nu}}
$$

then $x_{\rho}$ is a monomial of $Q(A)$ for some $\rho$. We let $M=M_{1} \times \cdots \times M_{\nu}$ with $\operatorname{dim}\left(M_{j}\right)=2 i_{j}$ and let $V=L_{1} \oplus \cdots \oplus L_{\nu} \oplus 1^{v-\nu}$ where the $L_{j}$ are line bundles over $M_{j}$. If $c\left(x_{\rho}, Q\right)$ is the coefficient of $x$ in $Q$, then:

$$
\int_{M} Q(\nabla)=c\left(x_{\rho}, Q\right) \prod_{j=1}^{\nu} \int_{M_{j}} c_{1}\left(L_{j}\right)^{i_{j}}
$$

This is, of course, nothing but an application of the splitting principle. This reduces the proof of this lemma to the special case $Q \in \mathcal{Q}_{m, m, m, 1}$.

If $Q=c_{1}^{k}$ we take $M=S^{2} \times \cdots \times S^{2}$ to be the $k$-fold product of two dimensional spheres. We let $L_{j}$ be a line bundle over the $j^{\text {th }}$ factor of $S^{2}$ and let $L=L_{1} \otimes \cdots \otimes L_{k}$. Then $c_{1}(L)=c_{1}\left(L_{j}\right)$ so

$$
\int_{M} c_{1}(L)^{k}=k!\prod_{j=1}^{k} \int_{S^{2}} c_{1}\left(L_{j}\right)
$$

which reduces the proof to the case $m=2$ and $k=1$. We gave an example in Lemma 2.1.5 of a line bundle over $S^{2}$ so $\int_{S^{2}} c_{1}(L)=1$. Alternatively, if we use the Gauss-Bonnet theorem with $L=T_{c}\left(S^{2}\right)$, then $\int_{S^{2}} c_{1}\left(T_{c}\left(S^{2}\right)\right)=\int_{S^{2}} e_{2}\left(T\left(S^{2}\right)\right)=\chi\left(S^{2}\right)=2 \neq 0$ completes the proof of (a). The proof of (b) is the same where we replace the real manifolds $M_{\rho}^{r}$ by the corresponding manifolds $M_{\rho}^{c}$ and the real basis $P_{\rho}$ by the corresponding basis $P_{\rho}^{c}$ of characteristic forms of $T_{c}(M)$. The remainder of the proof is the same and relies on Lemma 2.3.4 exactly as for (a) and is therefore omitted in the interests of brevity.

## CHAPTER 3 <br> THE INDEX THEOREM

## Introduction

In this the third chapter, we complete the proof of the index theorem for the four classical elliptic complexes. We give a proof of the Aityah-Singer theorem in general based on the Chern isomorphism between $K$-theory and cohomology (which is not proved). Our approach is to use the results of the first chapter to show there exists a suitable formula with the appropriate functorial properties. The results of the second chapter imply it must be a characteristic class. The normalizing constants are then determined using the method of universal examples.

In section 3.1, we define the twisted signature complex and prove the Hirzebruch signature theorem. We shall postpone until section 3.4 the determination of all the normalizing constants if we take coefficients in an auxilary bundle. In section 3.2, we introduce spinors as a means of connecting the de Rham, signature and Dolbeault complexes. In section 3.3, we discuss the obstruction to putting a spin structure on a real vector bundle in terms of Stieffel-Whitney classes. We compute the characteristic classes of spin bundles.

In section 3.4, we discuss the spin complex and the $\hat{A}$ genus. In section 3.5, we use the spin complex together with the $\operatorname{spin}_{c}$ representation to discuss the Dolbeault complex and to prove the Riemann-Roch theorem for almost complex manifolds. In sections 3.6 and 3.7 we give another treatment of the Riemann-Roch theorem based on a direct approach for Kaehler manifolds. For Kaehler manifolds, the integrands arising from the heat equation can be studied directly using an invariant characterization of the Chern forms similar to that obtained for the Euler form. These two subsections may be deleted by a reader not interested in Kaehler geometry.

In section 3.8, we give the preliminaries we shall need to prove the AtiyahSinger index theorem in general. The only technical tool we will use which we do not prove is the Chern isomorphism between rational cohomology and $K$-theory. We give a discussion of Bott periodicity using Clifford algebras. In section 3.9, we show that the index can be treated as a formula in rational $K$-theory. We use constructions based on Clifford algebras to determine the normalizing constants involved. For these two subsections, some familarity with $K$-theory is helpful, but not essential.

Theorems 3.1.4 and 3.6.10 were also derived by V. K. Patodi using a complicated cancellation argument as a replacement of the invariance theory presented in Chapter 2. A similar although less detailed discussion may also be found in the paper of Atiyah, Bott and Patodi.

### 3.1. The Hirzebruch Signature Formula.

The signature complex is best described using Clifford algebras as these provide a unified framework in which to discuss (and avoid) many of the $\pm$ signs which arise in dealing with the exterior algebra directly. The reader will note that we are choosing the opposite sign convention for our discussion of Clifford algebras from that adopted in the example of Lemma 2.1.5. This change in sign convention is caused by the $\sqrt{-1}$ present in discussing the symbol of a first order operator.

Let $V$ be a real vector with a positive definite inner product. The Clifford algebra $\operatorname{CLIF}(V)$ is the universal algebra generated by $V$ subject to the relations

$$
v * v+(v, v)=0 \text { for } v \in V .
$$

If the $\left\{e_{i}\right\}$ are an orthonormal basis for $V$ and if $I=\left\{1 \leq i_{1}<\cdots<i_{p} \leq\right.$ $\operatorname{dim}(V)\}$ then $e_{i} * e_{j}+e_{j} * e_{i}=-2 \delta_{i j}$ is the Kronecker symbol and

$$
e_{I}=e_{i_{1}} * \cdots * e_{i_{p}}
$$

is an element of the Clifford algebra. CLIF( $V$ ) inherits a natural inner product from $V$ and the $\left\{e_{I}\right\}$ form an orthonormal basis for $V$.

If $\Lambda(V)$ denotes the exterior algebra of $V$ and if $\operatorname{END}(\Lambda(V))$ is the algebra of linear endomorphisms of $\Lambda(V)$, there is a natural representation of $\operatorname{CLIF}(V)$ into $\operatorname{END}(\Lambda(V))$ given by Clifford multiplication. Let ext: $V \rightarrow \operatorname{END}(\Lambda(V))$ be exterior multiplication on the left and let $\operatorname{int}(v)$ be interior multiplication, the adjoint. For example:

$$
\begin{aligned}
& \operatorname{ext}\left(e_{1}\right)\left(e_{i_{1}} \wedge \cdots \wedge e_{i_{p}}\right)= \begin{cases}e_{1} \wedge e_{i_{1}} \wedge \cdots \wedge e_{i_{p}} & \text { if } i_{1}>1 \\
0 & \text { if } i_{1}=1\end{cases} \\
& \operatorname{int}\left(e_{1}\right)\left(e_{i_{1}} \wedge \cdots \wedge e_{i_{p}}\right)= \begin{cases}e_{i_{2}} \wedge \cdots \wedge e_{i_{p}} & \text { if } i_{1}=1 \\
0 & \text { if } i_{1}>1\end{cases}
\end{aligned}
$$

We define

$$
c(v)=\operatorname{ext}(v)-\operatorname{int}(v): V \rightarrow \operatorname{END}(\Lambda(V))
$$

It is immediate from the definition that

$$
c(v)^{2}=-(\operatorname{ext}(v) \operatorname{int}(v)+\operatorname{int}(v) \operatorname{ext}(v))=-|v|^{2} I
$$

so that $c$ extends to define an algebra morphism

$$
c: \operatorname{CLIF}(V) \rightarrow \operatorname{END}(\Lambda(V))
$$

Furthermore:

$$
c\left(e_{I}\right) 1=e_{i_{1}} \wedge \cdots \wedge e_{i_{p}}
$$

so the map $w \mapsto c(w) 1$ defines a vector space isomorphism (which is not, of course an algebra morphism) between $\operatorname{CLIF}(V)$ and $\Lambda(V)$. Relative to an orthonormal frame, we simply replace Clifford multiplication by exterior multiplication.

Since these constructions are all independent of the basis chosen, they extend to the case in which $V$ is a real vector bundle over $M$ with a fiber metric which is positive definite. We define $\operatorname{CLIF}(V), \Lambda(V)$, and $c: \operatorname{CLIF}(V) \rightarrow \operatorname{END}(\Lambda(V))$ as above. Since we only want to deal with complex bundles, we tensor with $\mathbf{C}$ at the end to enable us to view these bundles as being complex. We emphasize, however, that the underlying constructions are all real.

Clifford algebras provide a convenient way to describe both the de Rham and the signature complexes. Let $(d+\delta): C^{\infty}\left(\Lambda\left(T^{*} M\right)\right) \rightarrow C^{\infty}\left(\Lambda\left(T^{*} M\right)\right)$ be exterior differentiation plus its adjoint as discussed earlier. The leading symbol of $(d+\delta)$ is $\sqrt{-1}(\operatorname{ext}(\xi)-\operatorname{int}(\xi))=\sqrt{-1} c(\xi)$. We use the following diagram to define an operator $A$; let $\nabla$ be covariant differentiation. Then:

$$
A: C^{\infty}\left(\Lambda\left(T^{*} M\right)\right) \xrightarrow{\nabla} C^{\infty}\left(T^{*} M \otimes \Lambda\left(T^{*} M\right)\right) \xrightarrow{c} C^{\infty}\left(\Lambda\left(T^{*} M\right)\right)
$$

$A$ is invariantly defined. If $\left\{e_{i}\right\}$ is a local orthonormal frame for $T^{*}(M)$ which we identify with $T(M)$, then:

$$
A(\omega)=\sum_{i}\left(\operatorname{ext}\left(e_{i}\right)-\operatorname{int}\left(e_{i}\right)\right) \nabla_{e_{i}}(\omega)
$$

Since the leading symbol of $A$ is $\sqrt{-1} c(\xi)$, these two operators have the same leading symbol so $(d+\delta)-A=A_{0}$ is an invariantly defined $0^{\text {th }}$ order operator. Relative to a coordinate frame, we can express $A_{0}$ as a linear combination of the 1 -jets of the metric with coefficients which are smooth in the $\left\{g_{i j}\right\}$ variables. Given any point $x_{0}$, we can always choose a frame so the $g_{i j / k}$ variables vanish at $x_{0}$ so $A_{0}\left(x_{0}\right)=0$ so $A_{0} \equiv 0$. This proves $A=(d+\delta)$ is defined by this diagram which gives a convenient way of describing the operator $(d+\delta)$ in terms of Clifford multiplication.

This trick will be useful in what follows. If $A$ and $B$ are natural first order differential operators with the same leading symbol, then $A=B$ since $A-B$ is a $0^{\text {th }}$ order operator which is linear in the 1 -jets of the metric. This trick does not work in the holomorphic category unless we impose the additional hypothesis that $M$ is Kaehler. This makes the study of the Riemann-Roch theorem more complicated as we shall see later since there are many natural operators with the same leading symbol.

We let $\alpha \in \operatorname{END}\left(\Lambda\left(T^{*} M\right)\right)$ be defined by:

$$
\alpha\left(\omega_{p}\right)=(-1)^{p} \omega_{p} \quad \text { for } \omega_{p} \in \Lambda^{p}\left(T^{*} M\right)
$$

It is immediate that

$$
\operatorname{ext}(\xi) \alpha=-\alpha \operatorname{ext}(\xi) \quad \text { and } \quad \operatorname{int}(\xi) \alpha=-\alpha \operatorname{int}(\xi)
$$

so that $-\alpha(d+\delta) \alpha$ and $(d+\delta)$ have the same leading symbol. This implies

$$
\alpha(d+\delta)=-(d+\delta) \alpha
$$

We decompose $\Lambda\left(T^{*} M\right)=\Lambda^{\mathrm{e}}\left(T^{*} M\right) \oplus \Lambda^{\circ}\left(T^{*} M\right)$ into the differential forms of even and odd degree. This decomposes $\Lambda\left(T^{*} M\right)$ into the $\pm 1$ eigenspaces of $\alpha$. Since $(d+\delta)$ anti-commutes with $\alpha$, we decompose

$$
(d+\delta)_{\mathrm{e}, \mathrm{o}}: C^{\infty}\left(\Lambda^{\mathrm{e}, \mathrm{o}}\left(T^{*} M\right)\right) \rightarrow C^{\infty}\left(\Lambda^{\mathrm{e}, \mathrm{o}}\left(T^{*} M\right)\right)
$$

where the adjoint of $(d+\delta)_{\mathrm{e}}$ is $(d+\delta)_{\mathrm{o}}$. This is, of course, just the de Rham complex, and the index of this elliptic operator is $\chi(M)$.

If $\operatorname{dim}(M)=m$ is even and if $M$ is oriented, there is another natural endomorphism $\tau \in \operatorname{END}\left(\Lambda\left(T^{*} M\right)\right)$. It can be used to define an elliptic complex over $M$ called the signature complex in just the same way that the de Rham complex was defined. Let dvol $\in \Lambda^{m}\left(T^{*} M\right)$ be the volume form. If $\left\{e_{i}\right\}$ is an oriented local orthonormal frame for $T^{*} M$, then dvol $=$ $e_{1} \wedge \cdots \wedge e_{m}$. We can also regard dvol $=e_{1} * \cdots * e_{m} \in \operatorname{CLIF}\left(T^{*} M\right)$ and we define

$$
\tau=(\sqrt{-1})^{m / 2} c(\mathrm{dvol})=(\sqrt{-1})^{m / 2} c\left(e_{1}\right) \ldots c\left(e_{m}\right)
$$

We compute:

$$
\begin{aligned}
\tau^{2} & =(-1)^{m / 2} c\left(e_{1} * \cdots * e_{m} * e_{1} * \cdots * e_{m}\right) \\
& =(-1)^{m / 2}(-1)^{m+(m-1)+\cdots+1} \\
& =(-1)^{(m+(m+1) m) / 2}=1
\end{aligned}
$$

Because $m$ is even, $c(\xi) \tau=-\tau c(\xi)$ so $\tau$ anti-commutes with the symbol of $(d+\delta)$. If we decompose $\Lambda\left(T^{*} M\right)=\Lambda^{+}\left(T^{*} M\right) \oplus \Lambda^{-}\left(T^{*} M\right)$ into the $\pm 1$ eigenvalues of $\tau$, then $(d+\delta)$ decomposes to define:

$$
(d+\delta)_{ \pm}: C^{\infty}\left(\Lambda^{ \pm}\left(T^{*} M\right)\right) \rightarrow C^{\infty}\left(\Lambda^{\mp}\left(T^{*} M\right)\right)
$$

where the adjoint of $(d+\delta)_{+}$is $(d+\delta)_{-}$. We define:

$$
\operatorname{signature}(M)=\operatorname{index}(d+\delta)_{+}
$$

to be the signature of $M$. (This is also often refered to as the index of $M$, but we shall not use this notation as it might be a source of some confusion).

We decompose the Laplacian $\Delta=\Delta^{+} \oplus \Delta^{-}$so

$$
\operatorname{signature}(M)=\operatorname{dim} \mathrm{N}\left(\Delta^{+}\right)-\operatorname{dim} \mathrm{N}\left(\Delta^{-}\right)
$$

Let $a_{n}^{s}(x$, orn $)=a_{n}\left(x, \Delta^{+}\right)-a_{n}\left(x, \Delta^{-}\right)$be the invariants of the heat equation; they depend on the orientation orn chosen. Although we have complexified the bundles $\Lambda\left(T^{*} M\right)$ and $\operatorname{CLIF}\left(T^{*} M\right)$, the operator $(d+\delta)$ is real. If $m \equiv 2$ (4), then $\tau$ is pure imaginary. Complex conjugation defines an isomorphism

$$
\Lambda^{+}\left(T^{*} M\right) \stackrel{\simeq}{\rightarrow} \Lambda^{-}\left(T^{*} M\right) \quad \text { and } \quad \Delta^{+} \xrightarrow{\simeq} \Delta^{-} .
$$

This implies signature $(M)=0$ and $a_{n}^{s}(x$, orn $)=0$ in this case. We can get a non-zero index if $m \equiv 2$ (4) if we take coefficients in some auxiliary bundle as we shall discuss shortly.

If $m \equiv 0(4)$, then $\tau$ is a real endomorphism. In general, we compute:

$$
\begin{aligned}
\tau\left(e_{1} \wedge \cdots \wedge e_{p}\right) & =(\sqrt{-1})^{m / 2} e_{1} * \cdots * e_{m} * e_{1} * \cdots * e_{p} \\
& =(\sqrt{-1})^{m / 2}(-1)^{p(p-1) / 2} e_{p+1} \wedge \cdots \wedge e_{m} .
\end{aligned}
$$

If "*" is the Hodge star operator discussed in the first sections, then

$$
\tau_{p}=(\sqrt{-1})^{m / 2+p(p-1)} *_{p}
$$

acting on $p$-forms. The spaces $\Lambda^{p}\left(T^{*} M\right) \oplus \Lambda^{m-p}\left(T^{*} M\right)$ are invariant under $\tau$. If $p \neq m-p$ there is a natural isomorphism

$$
\Lambda^{p}\left(T^{*} M\right) \stackrel{\approx}{\rightarrow}\left(\Lambda^{p}\left(T^{*} M\right) \oplus \Lambda^{m-p}\left(T^{*} M\right)\right)^{ \pm} \quad \text { by } \quad \omega_{p} \mapsto \frac{1}{2}\left(\omega_{p} \pm \omega_{p}\right) .
$$

This induces a natural isomorphism from

$$
\Delta_{p} \stackrel{\widetilde{3}}{\rightarrow} \Delta^{ \pm} \text {on }\left(\Lambda^{p}\left(T^{*} M\right) \oplus \Lambda^{m-p}\left(T^{*} M\right)\right)^{ \pm}
$$

so these terms all cancel off in the alternating sum and the only contribution is made in the middle dimension $p=m-p$.

If $m=4 k$ and $p=2 k$ then $\tau=*$. We decompose $\mathrm{N}\left(\Delta_{p}\right)=\mathrm{N}\left(\Delta_{p}^{+}\right) \oplus$ $\mathrm{N}\left(\Delta_{p}^{-}\right)$so signature $(M)=\operatorname{dim} \mathrm{N}\left(\Delta_{p}^{+}\right)-\operatorname{dim} \mathrm{N}\left(\Delta_{p}^{-}\right)$. There is a natural symmetric bilinear form on $H^{2 k}\left(T^{*} M ; \mathbf{C}\right)=\mathrm{N}\left(\Delta_{p}\right)$ defined by

$$
I\left(\alpha_{1}, \alpha_{2}\right)=\int_{M} \alpha_{1} \wedge \alpha_{2}
$$

If we use the de Rham isomorphism to identify de Rham and simplicial cohomology, then this bilinear form is just the evaluation of the cup product
of two cohomology classes on the top dimensional cycle. This shows $I$ can be defined in purely topological terms.

The index of a real quadratic form is just the number of +1 eigenvalues minus the number of -1 eigenvalues when it is diagonalized over $\mathbf{R}$. Since

$$
I(\alpha, \beta)=\int_{M} \alpha \wedge \beta=\int_{M}(\alpha, * \beta) \mathrm{dvol}=(\alpha, * \beta)_{L^{2}}
$$

we see

$$
\begin{aligned}
\operatorname{index}(I)= & \operatorname{dim}\left\{+1 \text { eigenspace of } * \text { on } H^{p}\right\} \\
& \quad-\operatorname{dim}\left\{-1 \text { eigenspace of } * \text { on } H^{p}\right\} \\
= & \operatorname{dim} \mathrm{N}\left(\Delta^{+}\right)-\operatorname{dim} \mathrm{N}\left(\Delta^{-}\right)=\operatorname{signature}(M)
\end{aligned}
$$

This gives a purely topological definition of the signature of $M$ in terms of cup product. We note that if we reverse the orientation of $M$, then the signature changes sign.
Example: Let $M=\mathbf{C} P_{2 k}$ be complex projective space. Let $x \in H^{2}(M ; \mathbf{C})$ be the generator. Since $x^{k}$ is the generator of $H^{2 k}(M ; \mathbf{C})$ and since $x^{k} \wedge$ $x^{k}=x^{2 k}$ is the generator of $H^{4 k}(M ; \mathbf{C})$, we conclude that $* x^{k}=x^{k}$. $\operatorname{dim} \mathrm{N}\left(\Delta_{2 k}^{+}\right)=1$ and $\operatorname{dim} \mathrm{N}\left(\Delta_{2 k}^{-}\right)=0$ so signature $\left(\mathbf{C} P_{2 k}\right)=1$.

An important tool in the study of the de Rham complex was its multiplicative properties under products. Let $M_{i}$ be oriented even dimensional manifolds and let $M=M_{1} \times M_{2}$ with the induced orientation. Decompose:

$$
\begin{aligned}
\Lambda\left(T^{*} M\right) & =\Lambda\left(T^{*} M_{1}\right) \otimes \Lambda\left(T^{*} M_{2}\right) \\
\operatorname{CLIF}\left(T^{*} M\right) & =\operatorname{CLIF}\left(T^{*} M_{1}\right) \otimes \operatorname{CLIF}\left(T^{*} M_{2}\right)
\end{aligned}
$$

as graded non-commutative algebras-i.e.,

$$
\left(\omega_{1} \otimes \omega_{2}\right) \circ\left(\omega_{1}^{\prime} \otimes \omega_{2}^{\prime}\right)=(-1)^{\operatorname{deg} \omega_{2} \cdot \operatorname{deg} \omega_{1}^{\prime}}\left(\omega_{1} \circ \omega_{1}^{\prime}\right) \otimes\left(\omega_{2} \circ \omega_{2}^{\prime}\right)
$$

for $\circ=$ either $\wedge$ or $*$. Relative to this decomposition, we have:

$$
\tau=\tau_{1} \otimes \tau_{2} \text { where the } \tau_{i} \text { commute. }
$$

This implies that:

$$
\begin{aligned}
\Lambda^{+}(M) & =\Lambda^{+}\left(T^{*} M_{1}\right) \otimes \Lambda^{+}\left(T^{*} M_{2}\right) \oplus \Lambda^{-}\left(T^{*} M_{1}\right) \otimes \Lambda^{-}\left(T^{*} M_{2}\right) \\
\Lambda^{-}(M) & =\Lambda^{-}\left(T^{*} M_{1}\right) \otimes \Lambda^{+}\left(T^{*} M_{2}\right) \oplus \Lambda^{+}\left(T^{*} M_{1}\right) \otimes \Lambda^{-}\left(T^{*} M_{2}\right) \\
\mathrm{N}\left(\Delta^{+}\right) & =\mathrm{N}\left(\Delta_{1}^{+}\right) \otimes \mathrm{N}\left(\Delta_{2}^{+}\right) \oplus \mathrm{N}\left(\Delta_{1}^{-}\right) \otimes \mathrm{N}\left(\Delta_{2}^{-}\right) \\
\mathrm{N}\left(\Delta^{-}\right) & =\mathrm{N}\left(\Delta_{1}^{-}\right) \otimes \mathrm{N}\left(\Delta_{2}^{+}\right) \oplus \mathrm{N}\left(\Delta_{1}^{+}\right) \otimes \mathrm{N}\left(\Delta_{2}^{-}\right) \\
\operatorname{signature}(M) & =\operatorname{signature}\left(M_{1}\right) \operatorname{signature}\left(M_{2}\right)
\end{aligned}
$$

Example: Let $\rho$ be a partition of $k=i_{1}+\cdots+i_{j}$ and $M_{\rho}^{r}=\mathbf{C} P_{2 i_{1}} \times$ $\cdots \times \mathbf{C} P_{2 i_{j}}$. Then signature $\left(M_{\rho}^{r}\right)=1$. Therefore if $L_{k}$ is the Hirzebruch $L$-polynomial,

$$
\operatorname{signature}\left(M_{\rho}^{r}\right)=\int_{M_{\rho}^{p}} L_{k}\left(T\left(M_{\rho}^{r}\right)\right)
$$

by Lemma 2.3.5.
We can now begin the proof of the Hirzebruch signature theorem. We shall use the same argument as we used to prove the Gauss-Bonnet theorem with suitable modifications. Let $m=4 k$ and let $a_{n}^{s}(x$, orn $)=a_{n}\left(x, \Delta^{+}\right)-$ $a_{n}\left(x, \Delta^{-}\right)$be the invariants of the heat equation. By Lemma 1.7.6:

$$
\int_{M} a_{n}^{s}(x, \text { orn })= \begin{cases}0 & \text { if } n \neq m \\ \operatorname{signature}(M) & \text { if } n=m\end{cases}
$$

so this gives a local formula for signature ( $M$ ). We can express $\tau$ functorially in terms of the metric tensor. We can find functorial local frames for $\Lambda^{ \pm}$relative to any oriented coordinate system in terms of the coordinate frames for $\Lambda\left(T^{*} M\right)$. Relative to such a frame, we express the symbol of $\Delta^{ \pm}$functorially in terms of the metric. The leading symbol is $|\xi|^{2} I$; the first order symbol is linear in the 1-jets of the metric with coefficients which depend smoothly on the $\left\{g_{i j}\right\}$ variables; the $0^{\text {th }}$ order symbol is linear in the 2-jets of the metric and quadratic in the 1 -jets of the metric with coefficients which depend smoothly on the $\left\{g_{i j}\right\}$ variables. By Lemma 2.4.2, we conclude $a_{n}^{s}(x$, orn $)$ is homogeneous of order $n$ in the jets of the metric.

It is worth noting that if we replace the metric $G$ by $c^{2} G$ for $c>0$, then the spaces $\Lambda^{ \pm}$are not invariant. On $\Lambda^{p}$ we have:

$$
\tau\left(c^{2} G\right)\left(\omega_{p}\right)=c^{2 p-m} \tau(G)\left(\omega_{p}\right)
$$

However, in the middle dimension we have $\tau$ is invariant as are the spaces $\Lambda_{p}^{ \pm}$for $2 p=m$. Clearly $\Delta_{p}^{ \pm}\left(c^{2} G\right)=c^{-2} \Delta_{p}^{ \pm}(G)$. Since $a_{n}^{s}(x$, orn $)$ only depends on the middle dimension, this provides another proof that $a_{n}^{s}$ is homogeneous of order $n$ in the derivatives of the metric since

$$
\begin{aligned}
a_{n}^{s}(x, \text { orn })\left(c^{2} G\right) & =a_{n}\left(x, c^{-2} \Delta_{p}^{+}\right)-a_{n}\left(x, c^{-2} \Delta_{p}^{-}\right) \\
& =c^{-n} a_{n}\left(x, \Delta_{p}^{+}\right)-c^{-n} a_{n}\left(x, \Delta_{p}^{-}\right)=c^{-n} a_{n}^{s}(x, \text { orn })(G)
\end{aligned}
$$

If we reverse the orientation, we interchange the roles of $\Delta^{+}$and $\Delta^{-}$so $a_{n}^{s}$ changes sign if we reverse the orientation. This implies $a_{n}^{s}$ can be regarded as an invariantly defined $m$-form; $a_{n}^{s}(x)=a_{n}^{s}(x$, orn $) \mathrm{dvol} \in \mathcal{P}_{m, n, m}$.

Theorem 3.1.1. Let $a_{n}^{s}=\left\{a_{n}\left(x, \Delta^{+}\right)-a_{n}\left(x, \Delta^{-}\right)\right\}$dvol $\in \mathcal{P}_{m, n, m}$ then:
(a) $a_{n}^{s}=0$ if either $m \equiv 2$ (4) or if $n<m$.
(b) If $m=4 k$, then $a_{4 k}^{s}=L_{k}$ is the Hirzebruch polynomial so

$$
\operatorname{signature}(M)=\int_{M} L_{k}
$$

## (Hirzebruch Signature Theorem).

Proof: We already noted that $\Delta^{+}$is naturally isomorphic to $\Delta^{-}$if $m \equiv$ 2 (4) so $a_{n}^{s}=0$ in that case. Lemma 2.5.6 implies $a_{n}^{s}=0$ for $n<m$. If $m=4 k$, then $a_{m}^{s}$ is a characteristic form of $T(M)$ by Lemma 2.5.6. We know

$$
\int_{M_{\rho}^{r}} a_{m}^{s}=\operatorname{signature}\left(M_{\rho}^{r}\right)=1=\int_{M_{\rho}^{r}} L_{k}
$$

so Lemma 2.3.4 implies $a_{m}^{s}=L_{k}$. Since signature $(M)=\int_{M} a_{m}^{s}$ for any manifold $M$, we conclude signature $(M)=\int_{M} L_{k}$ in general which completes the proof of (b).

If $\omega \in \Lambda\left(T^{*} M\right)$, we define $\int_{M} \omega=\int_{M} \omega_{m}$ of the top degree form. With this notational convention, we can also express

$$
\operatorname{signature}(M)=\int_{M} L
$$

which is a common form in which the Hirzebruch signature theorem appears.

It is worth making a few remarks about the proof of this result. Just as in the case of the de Rham complex, the heat equation furnishes us with the a priori local formula for the signature of $M$. The invariance theory of the second chapter identifies this local formula as a characteristic class. We evaluate this local formula on a sufficient number of classifying examples to determine the normalizing constants to prove $a_{m}^{s}=L_{k}$.

There are a great many consequences of this theorem and of the GaussBonnet theorem. We present just a few to illustrate some of the applications:

Corrolary 3.1.2.
(a) Let $F \rightarrow M_{1} \rightarrow M_{2}$ be a finite covering projection. Then $\chi\left(M_{1}\right)=$ $\chi\left(M_{2}\right)|F|$. If $M_{2}$ is orientable, then $M_{1}$ is orientable and we give it the natural orientation inherited from $M_{2}$. Then signature $\left(M_{1}\right)=\operatorname{signature}\left(M_{2}\right)|F|$. (b) If $M_{1}$ and $M_{2}$ are manifolds of dimension $m$, we let $M_{1} \# M_{2}$ be the connected sum. This is defined by punching out disks in both manifolds and gluing along the common resulting boundaries. Then $\chi\left(M_{1} \# M_{2}\right)+$ $\chi\left(S^{m}\right)=\chi\left(M_{1}\right)+\chi\left(M_{2}\right)$. If $M_{1}$ and $M_{2}$ are oriented by some orientation
on $M_{1} \# M_{2}$ then $\operatorname{signature~}\left(M_{1} \# M_{2}\right)=\operatorname{signature}\left(M_{1}\right)+\operatorname{signature}\left(M_{2}\right)$, if $m \equiv 0$ (4).

Proof: (a) is an immediate consequence of the fact we have local formulas for the Euler characteristic and for the signature. To prove (b), we note that the two disks we are punching out glue together to form a sphere. We use the additivity of local formulas to prove the assertion about $X$. The second assertion follows similarly if we note signature $\left(S^{m}\right)=0$.

This corollary has topological consequences. Again, we present just one to illustrate the methods involved:

Corollary 3.1.3. Let $F \rightarrow \mathbf{C} P_{2 j} \rightarrow M$ be a finite covering. Then $|F|=1$ and $M=\mathbf{C} P_{2 j}$.
Proof: $\chi\left(\mathbf{C} P_{2 j}\right)=2 j+1$ so as $2 j+1=|F| \chi(M)$, we conclude $|F|$ must be odd. Therefore, this covering projection is orientation preserving so $M$ is orientable. The identity $1=\operatorname{signature}\left(\mathbf{C} P_{2 j}\right)=|F|$ signature $(M)$ implies $|F|=1$ and completes the proof.

If $m \equiv 2$ (4), the signature complex does not give a non-trivial index. We twist by taking coefficients in an auxiliary complex vector bundle $V$ to get a non-trivial index problem in any even dimension $m$.

Let $V$ be a smooth complex vector bundle of dimension $v$ equipped with a Riemannian connection $\nabla$. We take the Levi-Civita connection on $T^{*}(M)$ and on $\Lambda\left(T^{*} M\right)$ and let $\nabla$ be the tensor product connection on $\Lambda\left(T^{*} M\right) \otimes V$. We define the operator $(d+\delta)_{V}$ on $C^{\infty}\left(\Lambda\left(T^{*} M\right) \otimes V\right)$ using the diagram:

$$
\begin{aligned}
(d+\delta)_{V}: C^{\infty}\left(\Lambda\left(T^{*} M\right) \otimes V\right) \xrightarrow{\nabla} C\left(T^{*} M \otimes \Lambda( \right. & \left.\left(T^{*} M\right) \otimes V\right) \\
& \xrightarrow{c \otimes 1} C^{\infty}\left(\Lambda\left(T^{*} M\right) \otimes V\right) .
\end{aligned}
$$

We have already noted that if $V=1$ is the trivial bundle with flat connection, then the resulting operator is $(d+\delta)$.

We define $\tau_{V}=\tau \otimes 1$, then a similar argument to that given for the signature complex shows $\tau_{V}^{2}=1$ and $\tau_{V}$ anti-commutes with $(d+\delta)_{V}$. The $\pm 1$ eigenspaces of $\tau_{V}$ are $\Lambda^{ \pm}\left(T^{*} M\right) \otimes V$ and the twisted signature complex is defined by the diagram:

$$
(d+\delta)_{V}^{ \pm}: C^{\infty}\left(\Lambda^{ \pm}\left(T^{*} M\right) \otimes V\right) \rightarrow C^{\infty}\left(\Lambda^{\mp}\left(T^{*} M\right) \otimes V\right)
$$

where as before $(d+\delta)_{V}^{-}$is the adjoint of $(d+\delta)_{V}^{+}$. We let $\Delta_{V}^{ \pm}$be the associated Laplacians and define:

$$
\begin{aligned}
\operatorname{signature}(M, V) & =\operatorname{index}\left((d+\delta)_{V}^{+}\right)=\operatorname{dim} \mathrm{N}\left(\Delta_{V}^{+}\right)-\operatorname{dim} \mathrm{N}\left(\Delta_{V}^{-}\right) \\
a_{n}^{s}(x, V) & =\left\{a_{n}\left(x, \Delta_{V}^{+}\right)-a_{n}\left(x, \Delta_{V}^{-}\right)\right\} \operatorname{dvol} \in \Lambda^{m} \\
\int_{M} a_{n}^{s}(x, V) & =\operatorname{signature}(M, V)
\end{aligned}
$$

The invariance of the index under homotopies shows signature $(M, V)$ is independent of the metric on $M$, of the fiber metric on $V$, and of the Riemannian connection on $V$. If we do not choose a Riemannian connection on $V$, we can still compute signature $(M, V)=\operatorname{index}\left((d+\delta)_{V}^{+}\right)$, but then $(d+\delta)_{V}^{-}$is not the adjoint of $(d+\delta)_{V}^{+}$and $a_{n}^{s}$ is not an invariantly defined $m$ form.

Relative to a functorial coordinate frame for $\Lambda^{ \pm}$, the leading symbol of $\Delta^{ \pm}$is $|\xi|^{2} I$. The first order symbol is linear in the 1 -jets of the metric and the connection form on $V$. The $0^{\text {th }}$ order symbol is linear in the 2-jets of the metric and the connection form and quadratic in the 1-jets of the metric and the connection form. Thus $a_{n}^{s}(X, V) \in \mathcal{R}_{m, n, m, v}$. Theorem 2.6.1 implies $a_{n}^{s}=0$ for $n<m$ while $a_{m}^{s}$ is a characteristic form of $T(M)$ and of $V$.

If $m=2$ and if $v=1$, then $\mathcal{R}_{2,2,2,1}$ is one dimensional and is spanned by the first Chern class $c_{1}(V)=\operatorname{ch}(V)=\frac{i}{2 \pi} \Omega$. Consequently $a_{2}^{s}=c c_{1}$ in this case. We shall show later that this normalizing constant $c=2$-i.e.,

Lemma 3.1.4. Let $m=2$ and let $V$ be a line bundle over $M_{2}$. Then:

$$
\operatorname{signature}(M, V)=2 \int_{M} c_{1}(L)
$$

We postpone the proof of this lemma until later in this chapter.
With this normalizing constant established, we can compute a formula for signature $(M, V)$ in general:

Theorem 3.1.5. Let $L$ be the total L-polynomial and let $\operatorname{ch}(V)$ be the Chern character. Then:
(a) $a_{n}^{s}(x, V)=0$ for $n<m$.
(b) $a_{m}^{s}(x, V)=\sum_{4 s+2 t=m} L_{s}(T M) \wedge 2^{t} c h_{t}(V)$ so that:

$$
\operatorname{signature}(M, V)=\sum_{4 s+2 t=m} \int_{M} L_{s}(T M) \wedge 2^{t} c h_{t}(V)
$$

The factors of $2^{t}$ are perhaps a bit mysterious at this point. They arise from the normalizing constant of Lemma 3.1.4 and will be explained when we discuss the spin and Dolbeault complexes.
Proof: We have already proved (a). We know $a_{m}^{s}(x, V)$ is a characteristic form which integrates to signature $(M, V)$ so it suffices to verify the formula of (b). If $V_{1}$ and $V_{2}$ are bundles, we let $V=V_{1} \oplus V_{2}$ with the direct sum connection. Since $\Delta_{V}^{ \pm}=\Delta_{V_{1}}^{ \pm} \oplus \Delta_{V_{2}}^{ \pm}$we conclude

$$
\operatorname{signature}\left(M, V_{1} \oplus V_{2}\right)=\operatorname{signature}\left(M, V_{1}\right)+\operatorname{signature}\left(M, V_{2}\right)
$$

Since the integrals are additive, we apply the uniqueness of Lemma 2.6.3 to conclude the local formulas must be additive. This also follows from Lemma 1.7.5 so that:

$$
a_{n}^{s}\left(x, V_{1} \oplus V_{2}\right)=a_{n}^{s}\left(x, V_{1}\right)+a_{n}^{s}\left(x, V_{2}\right) .
$$

Let $\left\{P_{\rho}\right\}_{|\rho|=s}$ be the basis for $\mathcal{P}_{m, 4 s, 4 s}$ and expand:

$$
a_{m}^{s}(x, V)=\sum_{4|\rho|+2 t=m} P_{\rho} \wedge Q_{m, t, v, \rho} \quad \text { for } \quad Q_{m, t, v, \rho} \in \mathcal{Q}_{m, 2 t, 2 t, v}
$$

a characteristic form of $V$. Then the additivity under direct sum implies:

$$
Q_{m, t, v, \rho}\left(V_{1} \oplus V_{2}\right)=Q_{m, t, v_{1}, \rho}\left(V_{1}\right)+Q_{m, t, v_{2}, \rho}\left(V_{2}\right)
$$

If $v=1$, then $Q_{m, t, 1, \rho}\left(V_{1}\right)=c \cdot c_{1}(V)^{t}$ since $\mathcal{Q}_{m, 2 t, 2 t, 1}$ is one dimensional. If $A$ is diagonal matrix, then the additivity implies:

$$
Q_{m, t, v, \rho}(A)=Q_{m, t, v, \rho}(\lambda)=c \cdot \sum_{j} \lambda_{j}^{t}=c \cdot c h_{t}(A) .
$$

Since $Q$ is determined by its values on diagonal matrices, we conclude:

$$
Q_{m, t, v, \rho}(V)=c(m, t, \rho) c h_{t}(V)
$$

where the normalizing constant does not depend on the dimension $v$. Therefore, we expand $a_{m}^{s}$ in terms of $c h_{t}(V)$ to express:

$$
a_{m}^{s}(x, V)=\sum_{4 s+2 t=m} P_{m, s} \wedge 2^{t} c h_{t}(V) \quad \text { for } P_{m, s} \in \mathcal{P}_{m, 4 s, 4 s}
$$

We complete the proof of the theorem by identifying $P_{m, s}=L_{s}$; we have reduced the proof of the theorem to the case $v=1$.

We proceed by induction on $m$; Lemma 3.1.4 establishes this theorem if $m=2$. Suppose $m \equiv 0(4)$. If we take $V$ to be the trivial bundle, then if $4 k=m$,

$$
a_{m}^{s}(x, 1)=L_{k}=P_{m, k}
$$

follows from Theorem 3.1.1. We may therefore assume $4 s<m$ in computing $P_{m, s}$. Let $M=M_{1} \times S^{2}$ and let $V=V_{1} \otimes V_{2}$ where $V_{1}$ is a line bundle over $M_{1}$ and where $V_{2}$ is a line bundle over $S^{2}$ so $\int_{S^{2}} c_{1}\left(V_{2}\right)=1$. (We constructed such a line bundle in section 2.1 using Clifford matrices). We take the product connection on $V_{1} \otimes V_{2}$ and decompose:

$$
\begin{aligned}
& \Lambda^{+}(V)=\Lambda^{+}\left(V_{1}\right) \otimes \Lambda^{+}\left(V_{2}\right) \oplus \Lambda^{-}\left(V_{1}\right) \otimes \Lambda^{-}\left(V_{2}\right) \\
& \Lambda^{-}(V)=\Lambda^{-}\left(V_{1}\right) \otimes \Lambda^{+}\left(V_{2}\right) \oplus \Lambda^{+}\left(V_{1}\right) \otimes \Lambda^{-}\left(V_{2}\right)
\end{aligned}
$$

A similar decomposition of the Laplacians yields:

$$
\begin{aligned}
\operatorname{signature}(M, V) & =\operatorname{signature}\left(M_{1}, V_{1}\right) \operatorname{signature}\left(M_{2}, V_{2}\right) \\
& =2 \operatorname{signature}\left(M_{1}, V_{1}\right)
\end{aligned}
$$

by Lemma 3.1.4. Since the signatures are multiplicative, the local formulas are multiplicative by the uniqueness assertion of Lemma 2.6.3. This also follows by using Lemma 1.7.5 that

$$
a_{m}^{s}(x, V)=\sum_{p+q=m} a_{p}^{s}\left(x_{1}, V_{1}\right) a_{q}^{s}\left(x_{2}, V_{2}\right)
$$

and the fact $a_{p}=0$ for $p<m_{1}$ and $a_{q}=0$ for $q<m_{2}$. Thus we conclude:

$$
a_{m}^{s}(x, V)=a_{m_{1}}^{s}\left(x_{1}, V_{1}\right) a_{m_{2}}^{s}\left(x_{2}, V_{2}\right)
$$

where, of course, $m_{2}=2$ and $m_{1}=m-2$.
We use the identity:

$$
\operatorname{ch}\left(V_{1} \otimes V_{2}\right)=\operatorname{ch}\left(V_{1}\right) \operatorname{ch}\left(V_{2}\right)
$$

to conclude therefore:

$$
\begin{aligned}
\operatorname{signature}\left(M_{1},\right. & \left.V_{1}\right) \\
& =\frac{1}{2} \operatorname{signature}(M, V) \\
& =\frac{1}{2}\left\{\sum_{4 s+2 t=m-2} \int_{M_{I}} P_{m, s} \wedge 2^{t} c h_{t}\left(V_{1}\right)\right\} \int_{M_{2}} 2 c h_{1}\left(V_{2}\right) \\
& =\sum_{4 s+2 t=m-2} \int_{M_{1}} P_{m, s} \wedge 2^{t} c h_{t}\left(V_{1}\right)
\end{aligned}
$$

We apply the uniqueness assertion of Lemma 2.6.3 to conclude $P_{m, s}=$ $P_{m-2, s}$ for $4 s \leq m-2$. Since by induction, $P_{m-2, s}=L_{s}$ this completes the proof of the theorem.

We note that the formula is non-zero, so Lemma 2.6.3 implies that in any even dimension $m$, there always exist $(M, V)$ so signature $(M, V) \neq 0$. In fact, much more is true. Given any orientable manifold $M$, we can find $V$ over $M$ so signature $(M, V) \neq 0$ if $\operatorname{dim}(M)$ is even. Since the proof of this assertion relies on the fact $c h: K(M) \otimes \mathbf{Q} \xrightarrow{\simeq} H^{2 *}(M ; \mathbf{Q})$ we postpone a discussion of this fact until we discuss the index theorem in general.

### 3.2. Spinors and their Representations.

To define the signature complex, we needed to orient the manifold $M$. For any Riemannian manifold, by restricting to local orthonormal frames, we can always assume the transition functions of $T(M)$ are maps $g_{\alpha \beta}: U_{\alpha} \cap$ $U_{\beta} \rightarrow \mathrm{O}(m)$. If $M$ is oriented, by restricting to local orthonormal frames, we can choose the transition functions of $T(M)$ to lie in $\mathrm{SO}(m)$ and to reduce the structure group from $\mathrm{GL}(m, R)$ to $\mathrm{SO}(m)$. The signature complex results from the represntation $\Lambda^{ \pm}$of $\mathrm{SO}(m)$; it cannot be defined in terms of $\mathrm{GL}(m, R)$ or $\mathrm{O}(m)$. By contrast, the de Rham complex results from the representation $\Lambda^{\mathrm{e}, \mathrm{o}}$ which is a representation of GL $(m, R)$ so the de Rham complex is defined for non-orientable manifolds.

To define the spin complex, which is in some sense a more fundamental elliptic complex than is either the de Rham or signature complex, we will have to lift the transition functions from $\operatorname{SO}(m)$ to $\operatorname{SPIN}(m)$. Just as every manifold is not orientable, in a similar fashion there is an obstruction to defining a spin structure.

If $m \geq 3$, then $\pi_{1}(\mathrm{SO}(m))=\mathbf{Z}_{2}$. Abstractly, we define $\operatorname{SPIN}(m)$ to be the universal cover of $\mathrm{SO}(m)$. (If $m=2$, then $\mathrm{SO}(2)=S^{1}$ and we let $\operatorname{SPIN}(2)=S^{1}$ with the natural double cover $\mathbf{Z}_{2} \rightarrow \mathrm{SPIN}(2) \rightarrow \mathrm{SO}(2)$ given by $\theta \mapsto 2 \theta)$. To discuss the representations of $\operatorname{SPIN}(m)$, it is convenient to obtain a more concrete representation of $\operatorname{SPIN}(m)$ in terms of Clifford algebras.

Let $V$ be a real vector space of dimension $v \equiv 0$ (2). Let

$$
\otimes V=\mathbf{R} \oplus V \oplus(V \otimes V) \oplus \cdots \oplus \bigotimes^{k} V \oplus \cdots
$$

be the complete tensor algebra of $V$. We assume $V$ is equipped with a symmetric positive definite bilinear form. Let $I$ be the two-sided ideal of $\otimes V$ generated by $\left\{v \otimes v+|v|^{2}\right\}_{v \in V}$, then the real Clifford algebra $\operatorname{CLIF}(V)=\bigotimes V \bmod I$. (Of course, we will always construct the corresponding complex Clifford algebra by tensoring CLIF( $V$ ) with the complex numbers). There is a natural transpose defined on $\otimes V$ by:

$$
\left(v_{1} \otimes \cdots \otimes v_{k}\right)^{t}=v_{k} \otimes \cdots \otimes v_{1} .
$$

Since this preserves the ideal $I$, it extends to $\operatorname{CLIF}(V)$. If $\left\{e_{1}, \ldots, e_{v}\right\}$ are the orthonormal basis for $V$, then $e_{i} * e_{j}+e_{j} * e_{i}=-2 \delta_{i j}$ and

$$
\left(e_{i_{1}} * \cdots * e_{i_{p}}\right)^{t}=(-1)^{p(p-1) / 2} e_{i_{1}} * \cdots * e_{i_{p}}
$$

If $V$ is oriented, we let $\left\{e_{i}\right\}$ be an oriented orthonormal basis and define:

$$
\tau=(\sqrt{-1})^{v / 2} e_{1} * \cdots * e_{v}
$$

We already computed that:

$$
\begin{aligned}
\tau^{2} & =(-1)^{v / 2} e_{1} * \cdots * e_{v} * e_{1} * \ldots e_{v} \\
& =(-1)^{v(v-1) / 2} e_{1} * \cdots * e_{v} * e_{1} * \cdots * e_{v} \\
& =e_{1} * \cdots * e_{v} * e_{v} * \cdots * e_{1}=(-1)^{v}=1
\end{aligned}
$$

We let $\operatorname{SPIN}(V)$ be the set of all $w \in \operatorname{CLIF}(V)$ such that $w$ can be decomposed as a formal product $w=v_{1} * \cdots * v_{2 j}$ for some $j$ where the $v_{i} \in V$ are elements of length 1 . It is clear that $w^{t}$ is such an element and that $w w^{t}=1$ so that $\operatorname{SPIN}(V)$ forms a group under Clifford multiplication. We define

$$
\rho(w) x=w x w^{t} \quad \text { for } x \in \mathrm{CLIF}, w \in \operatorname{CLIF}(V)
$$

For example, if $v_{1}=e_{1}$ is the first element of our orthonormal basis, then:

$$
e_{1} e_{i} e_{1}= \begin{cases}-e_{1} & i=1 \\ e_{i} & i \neq 1\end{cases}
$$

The natural inclusion of $V$ in $\otimes V$ induces an inclusion of $V$ in $\operatorname{CLIF}(V)$. $\rho\left(e_{1}\right)$ preserves $V$ and is reflection in the hyperplane defined by $e_{1}$. If $w \in$ $\operatorname{SPIN}(V)$, then $\rho(w): V \rightarrow V$ is a product of an even number of hyperplane reflections. It is therefore in $\mathrm{SO}(V)$ so

$$
\rho: \operatorname{SPIN}(V) \rightarrow \mathrm{SO}(V)
$$

is a group homomorphism. Since any orthogonal transformation of determinant one can be decomposed as a product of an even number of hyperplane reflections, $\rho$ is subjective.

If $w \in \operatorname{CLIF}(V)$ is such that

$$
w v w^{t}=v \text { all } v \in V \text { and } w w^{t}=1
$$

then $w v=v w$ for all $v \in V$ so $w$ must be in the center of $\operatorname{CLIF}(V)$.
Lemma 3.2.1. If $\operatorname{dim} V=v$ is even, then the center of $\operatorname{CLIF}(V)$ is one dimensional and consists of the scalars.

Proof: Let $\left\{e_{i}\right\}$ be an orthonormal basis for $V$ and let $\left\{e_{I}\right\}$ be the corresponding orthonormal basis for $\operatorname{CLIF}(V)$. We compute:

$$
e_{i} * e_{i_{1}} * \cdots * e_{i_{p}} * e_{i}=e_{i_{1}} * \cdots * e_{i_{p}} \quad \text { for } p \geq 1
$$

if (a) $p$ is even and $i$ is one of the $i_{j}$ or (b) $p$ is odd and $i$ is not one of the $i_{j}$. Thus given $I$ we can choose $i$ so $e_{i} * e_{I}=-e_{I} * e_{i}$. Thus
$e_{i} *\left(\sum_{I} c_{I} e_{I}\right)=\sum_{I} c_{I} e_{I} * e_{i}$ for all $i$ implies $c_{I}=0$ for $|I|>0$ which completes the proof.
(We note that this lemma fails if $\operatorname{dim} V$ is odd since the center in that case consists of the elements $a+b e_{1} * \cdots * e_{v}$ and the center is two dimensional).

If $\rho(w)=1$ and $w \in \operatorname{SPIN}(V)$, this implies $w$ is scalar so $w= \pm 1$. By considering the arc in $\operatorname{SPIN}(V)$ given by $w(\theta)=\left((\cos \theta) e_{1}+(\sin \theta) e_{2}\right) *$ $\left(-(\cos \theta) e_{1}+(\sin \theta) e_{2}\right)$, we note $w(0)=1$ and $w\left(\frac{\pi}{2}\right)=-1$ so $\operatorname{SPIN}(V)$ is connected. This proves we have an exact sequence of groups in the form:

$$
\mathbf{Z}_{2} \rightarrow \operatorname{SPIN}(V) \rightarrow \mathrm{SO}(V)
$$

and shows that $\operatorname{SPIN}(V)$ is the universal cover of $\mathrm{SO}(V)$ for $v>2$.
We note that is it possible to define $\operatorname{SPIN}(V)$ using Clifford algebras even if $v$ is odd. Since the center of $\operatorname{CLIF}(V)$ is two dimensional for $v$ odd, more care must be used with the relevant signs which arise. As we shall not need that case, we refer to Atiyah-Bott-Shapiro (see bibliography) for further details.
$\operatorname{SPIN}(V)$ acts on $\operatorname{CLIF}(V)$ from the left. This is an orthogonal action.
Lemma 3.2.2. Let $\operatorname{dim} V=2 v_{1}$. We complexify $\operatorname{CLIF}(V)$. As a left $\operatorname{SPIN}(V)$ module, this is not irreducible. We can decompose $\operatorname{CLIF}(V)$ as a direct sum of left $\operatorname{SPIN}(V)$ modules in the form

$$
\operatorname{CLIF}(V)=2^{v_{1}} \Delta .
$$

This representation is called the spin representation. It is not irreducible but further decomposes in the form

$$
\Delta=\Delta^{+} \oplus \Delta^{-}
$$

If we orient $V$ and let $\tau$ be the orientation form discussed earlier, then left multiplication by $\tau$ is $\pm 1$ on $\Delta^{ \pm}$so these are inequivalent representations. They are irreducible and act on a representation space of dimension $2^{v_{1}-1}$ and are called the half-spin representations.

Proof: Fix an oriented orthonormal basis $\left\{e_{i}\right\}$ for $V$ and define:

$$
\alpha_{1}=\sqrt{-1} e_{1} e_{2}, \quad \alpha_{2}=\sqrt{-1} e_{3} e_{4}, \quad \ldots, \quad \alpha_{v_{1}}=\sqrt{-1} e_{v-1} e_{v}
$$

as elements of $\operatorname{CLIF}(V)$. It is immediate that $\tau=\alpha_{1} \ldots \alpha_{v_{1}}$ and:

$$
\alpha_{i}^{2}=1 \quad \text { and } \quad \alpha_{i} \alpha_{j}=\alpha_{j} \alpha_{i}
$$

We let the $\left\{\alpha_{i}\right\}$ act on $\operatorname{CLIF}(V)$ from the right and decompose $\operatorname{CLIF}(V)$ into the $2^{v_{1}}$ simultaneous eigenspaces of this action. Since right and left multiplication commute, each eigenspace is invariant as a left $\operatorname{SPIN}(V)$
module. Each eigenspace corresponds to one of the $2^{v_{1}}$ possible sequences of + and - signs. Let $\varepsilon$ be such a string and let $\Delta_{\varepsilon}$ be the corresponding representation space.

Since $e_{1} \alpha_{1}=-\alpha_{1} e_{1}$ and $e_{1} \alpha_{i}=\alpha_{i} e_{1}$ for $i>1$, multiplication on the right by $e_{1}$ transforms $\Delta_{\varepsilon} \rightarrow \Delta_{\varepsilon^{\prime}}$, where $\varepsilon(1)=-\varepsilon^{\prime}(1)$ and $\varepsilon(i)=\varepsilon^{\prime}(i)$ for $i>1$. Since this map commutes with left multiplication by $\operatorname{SPIN}(V)$, we see these two representations are isomorphic. A similar argument on the other indices shows all the representation spaces $\Delta_{\varepsilon}$ are equivalent so we may decompose $\operatorname{CLIF}(V)=2^{v_{1}} \Delta$ as a direct sum of $2^{v_{1}}$ equivalent representation spaces $\Delta$. Since $\operatorname{dim}(C L I F(V))=2^{v}=4^{v_{1}}$, the representation space $\Delta$ has dimension $2^{v_{1}}$. We decompose $\Delta$ into $\pm 1$ eigenspaces under the action of $\tau$ to define $\Delta^{ \pm}$. Since $e_{1} * \cdots * e_{v}$ is in the center of $\operatorname{SPIN}(V)$, these spaces are invariant under the left action of $\operatorname{SPIN}(V)$. We note that elements of $V$ anti-commute with $\tau$ so Clifford multiplication on the left defines a map:

$$
c l: V \otimes \Delta^{ \pm} \rightarrow \Delta^{\mp}
$$

so both of the half-spin representations have the same dimension. They are clearly inequivalent. We leave the proof that they are irreducible to the reader as we shall not need this fact.

We shall use the notation $\Delta^{ \pm}$to denote both the representations and the corresponding representation spaces. Form the construction we gave, it is clear the CLIF $(V)$ acts on the left to preserve the space $\Delta$ so $\Delta$ is a representation space for the left action by the whole Clifford algebra. Since $v \tau=-\tau v$ for $v \in V$, Clifford multiplication on the left by an element of $V$ interchanges $\Delta^{+}$and $\Delta^{-}$.
Lemma 3.2.3. There is a natural map given by Clifford multiplication of $V \otimes \Delta^{ \pm} \rightarrow \Delta^{\mp}$. This map induces a map on the representations involved: $\rho \otimes \Delta^{ \pm} \mapsto \Delta^{\mp}$. If this map is denoted by $v * w$ then $v * v * w=-|v|^{2} w$.
Proof: We already checked the map on the spaces. We check:

$$
w v w^{t} \otimes w x \mapsto w v w^{t} w x=w v x
$$

to see that the map preserves the relevant representation. We emphasize that $\Delta^{ \pm}$are complex representations since we must complexify CLIF ( $V$ ) to define these representation spaces $\left(w w^{t}=1\right.$ for $\left.w \in \operatorname{SPIN}\right)$.

There is a natural map $\rho: \operatorname{SPIN}(V) \rightarrow \mathrm{SO}(V)$. There are natural representations $\Lambda, \Lambda^{ \pm}, \Lambda^{\mathrm{e}, \mathrm{o}}$ of $\mathrm{SO}(V)$ on the subspaces of $\Lambda(V)=\operatorname{CLIF}(V)$. We use $\rho$ to extend these representations of $\operatorname{SPIN}(V)$ as well. They are related to the half-spin representations as follows:

Lemma 3.2.4.
(a) $\Lambda=\Delta \otimes \Delta$.
(b) $\left(\Lambda^{+}-\Lambda^{-}\right)=\left(\Delta^{+}-\Delta^{-}\right) \otimes\left(\Delta^{+}+\Delta^{-}\right)$.
(c) $\left(\Lambda^{e}-\Lambda^{\circ}\right)=\left(\Delta^{+}-\Delta^{-}\right) \otimes\left(\Delta^{+}-\Delta^{-}\right)(-1)^{v / 2}$.

These identities are to be understood formally (in the sense of $K$-theory). They are shorthand for the identities:

$$
\Lambda^{+}=\left(\Delta^{+} \otimes \Delta^{+}\right) \oplus\left(\Delta^{+} \otimes \Delta^{-}\right) \quad \text { and } \quad \Delta^{-}=\left(\Delta^{-} \otimes \Delta^{+}\right) \oplus\left(\Delta^{-} \otimes \Delta^{-}\right)
$$

and so forth.
Proof: If we let $\operatorname{SPIN}(V)$ act on $\operatorname{CLIF}(V)$ by right multiplication by $w^{t}=w^{-1}$, then this representation is equivalent to left multiplication so we get the same decomposition of $\operatorname{CLIF}(V)=2^{v_{1}} \Delta$ as right spin modules. Since $\rho(w) x=w x w^{t}$, it is clear that $\Lambda=\Delta \otimes \Delta$ where one factor is viewed as a left and the other as a right spin module. Since $\Lambda^{ \pm}$is the decomposition of $\operatorname{SPIN}(V)$ under the action of $\tau$ from the left, (b) is immediate. (c) follows similarly once the appropriate signs are taken into consideration; $w \in \operatorname{CLIF}(V)^{\text {even }}$ if and only if

$$
\tau w \tau^{t}=(-1)^{v / 2} w
$$

which proves (c).
It is helpful to illustrate this for the case $\operatorname{dim} V=2$. Let $\left\{e_{1}, e_{2}\right\}$ be an oriented orthonormal basis for $V$. We compute that:

$$
\begin{aligned}
\left((\cos \alpha) e_{1}+(\sin \alpha) e_{2}\right)\left((\cos \beta) e_{1}\right. & \left.+(\sin \beta) e_{2}\right) \\
= & (-\cos \alpha \cos \beta-\sin \alpha \sin \beta) \\
& \quad+(\cos \alpha \sin \beta-\sin \alpha \cos \beta) e_{1} e_{2}
\end{aligned}
$$

so elements of the form $\cos \gamma+(\sin \gamma) e_{1} e_{2}$ belong to $\operatorname{SPIN}(V)$. We compute:

$$
\left(\cos \alpha+(\sin \alpha) e_{1} e_{2}\right)\left(\cos \beta+(\sin \beta) e_{1} e_{2}\right)=\cos (\alpha+\beta)+\sin (\alpha+\beta) e_{1} e_{2}
$$

so $\operatorname{spin}(V)$ is the set of all elements of this form and is naturally isomorphic to the circle $S^{1}=[0,2 \pi]$ with the endpoints identified. We compute that

$$
\begin{aligned}
\rho(w)\left(e_{1}\right) & =\left(\cos \theta+(\sin \theta) e_{1} e_{2}\right) e_{1}\left(\cos \theta+(\sin \theta) e_{2} e_{1}\right) \\
& =\left(\cos ^{2} \theta-\sin ^{2} \theta\right) e_{1}+2(\cos \theta \sin \theta) e_{2} \\
& =\cos (2 \theta) e_{1}+\sin (2 \theta) e_{2} \\
\rho(w)\left(e_{2}\right) & =\cos (2 \theta) e_{2}-\sin (2 \theta) e_{1}
\end{aligned}
$$

so the map $\rho: S^{1} \rightarrow S^{1}$ is the double cover $\theta \mapsto 2 \theta$.
We construct the one dimensional subspaces $V_{i}$ generated by the elements:

$$
\left.\begin{array}{rlrlrl}
v_{1} & =1+\tau, & v_{2} & =1-\tau, & v_{3} & =(1+\tau) e_{1}, \\
\tau v_{1} & =v_{1}, & \tau v_{2} & =-v_{2}, & \tau v_{3} & =v_{3}, \\
v_{1} \tau & =v_{1}, & & \tau v_{4} & =-v_{4}, \\
v_{2} \tau & =-v_{2}, & v_{3} \tau & =-v_{3}, & & v_{4} \tau
\end{array}\right)=v_{4} .
$$

When we decompose $\operatorname{CLIF}(V)$ under the left action of $\operatorname{SPIN}(V)$,

$$
V_{1} \simeq V_{3} \simeq \Delta^{+} \quad \text { and } \quad V_{2} \simeq V_{4} \simeq \Delta^{-}
$$

When we decompose $\operatorname{CLIF}(V)$ under the right action of $\operatorname{SPIN}(V)$, and replace $\tau$ by $-\tau=\tau^{t}$ acting on the right:

$$
V_{2} \simeq V_{3} \simeq \Delta^{+} \quad \text { and } \quad V_{1} \simeq V_{4} \simeq \Delta^{-}
$$

From this it follows that as $\mathrm{SO}(V)$ modules we have:

$$
V_{1}=\Delta^{+} \otimes \Delta^{-}, \quad V_{2}=\Delta^{-} \otimes \Delta^{+}, \quad V_{3}=\Delta^{+} \otimes \Delta^{+}, \quad V_{4}=\Delta^{-} \otimes \Delta^{-}
$$

from which it is immediate that:

$$
\begin{aligned}
\Lambda^{\mathrm{e}} & =V_{1} \oplus V_{2}
\end{aligned}=\Delta^{+} \otimes \Delta^{-} \oplus \Delta^{-} \otimes \Delta^{+} .
$$

If $V$ is a one-dimensional complex vector space, we let $V_{r}$ be the underlying real vector space. This defines a natural inclusion of $\mathrm{U}(1) \rightarrow \mathrm{SO}(2)$. If we let $J: V \rightarrow V$ be complex multiplication by $\sqrt{-1}$ then $J\left(e_{1}\right)=e_{2}$ and $J\left(e_{2}\right)=-e_{1}$ so $J$ is equivalent to Clifford multiplication by $e_{1} e_{2}$ on the left. We define a complex linear map from $V$ to $\Delta^{-} \otimes \Delta^{-}$by:

$$
T(v)=v-i J(v) \in V_{4}
$$

by computing $T\left(e_{1}\right)=\left(e_{1}-i e_{2}\right)=\left(1-i e_{1} e_{2}\right)\left(e_{1}\right)$ and $T\left(e_{2}\right)=e_{2}+i e_{1}=$ (i) $\left(e_{1}-i e_{2}\right)$.

Lemma 3.2.5. Let $\mathrm{U}(1)$ be identified with $\mathrm{SO}(2)$ in the usual manner. If $V$ is the underlying complex 1-dimensional space corresponding to $V_{r}$ then $V \simeq \Delta^{-} \otimes \Delta^{-}$and $V^{*} \simeq \Delta^{+} \otimes \Delta^{+}$as representation spaces of $\operatorname{SPIN}(2)$.

Proof: We have already verified the first assertion. The second follows from the fact that $V^{*}$ was made into a complex space using the map $-J$ instead of $J$ on $V_{r}$. This takes us into $V_{3}$ instead of into $V_{4}$. It is also clear that $\Delta^{-}=\left(\Delta^{+}\right)^{*}$ if $\operatorname{dim} V=2$. Of course, all these statements are to be interpreted as statements about representations since they are trivial as statements about vector spaces (since any vector spaces of the same dimension are isomorphic).

### 3.3. Spin Structures on Vector Bundles.

We wish to apply the constructions of 3.2 to vector bundles over manifolds. We first review some facts regarding principal bundles and StieffelWhitney classes which we shall need.

Principal bundles are an extremely convenient bookkeeping device. If $G$ is a Lie group, a principal $G$-bundle is a fiber space $\pi: P_{G} \rightarrow M$ with fiber $G$ such that the transition functions are elements of $G$ acting on $G$ by left multiplication in the group. Since left and right multiplication commute, we can define a right action of $G$ on $P_{G}$ which is fiber preserving. For example, let $\operatorname{SO}(2 k)$ and $\operatorname{SPIN}(2 k)$ be the groups defined by $\mathbf{R}^{2 k}$ with the cannonical inner product. Let $V$ be an oriented Riemannian vector bundle of dimension $2 k$ over $M$ and let $P_{\text {SO }}$ be the bundle of oriented frames of $V . \quad P_{\mathrm{SO}}$ is an $\mathrm{SO}(2 k)$ bundle and the natural action of $\mathrm{SO}(2 k)$ from the right which sends an oriented orthonormal frame $s=\left(s_{1}, \ldots, s_{2 k}\right)$ to $s \cdot g=\left(s_{1}^{\prime}, \ldots, s_{2 k}^{\prime}\right)$ is defined by:

$$
s_{i}^{\prime}=s_{1} g_{1, i}+\cdots+s_{2 k} g_{2 k, i}
$$

The fiber of $P_{\mathrm{SO}}$ is $\mathrm{SO}\left(V_{x}\right)$ where $V_{x}$ is the fiber of $V$ over the point x. This isomorphism is not natural but depends upon the choice of a basis.

It is possible to define the theory of characteristic classes using principal bundles rather than vector bundles. In this approach, a connection is a splitting of $T\left(P_{G}\right)$ into vertical and horizontal subspaces in an equivariant fashion. The curvature becomes a Lie algebra valued endomorphism of $T\left(P_{G}\right)$ which is equivariant under the right action of the group. We refer to Euguchi, Gilkey, Hanson for further details.

Let $\left\{U_{\alpha}\right\}$ be a cover of $M$ so $V$ is trivial over $U_{\alpha}$ and let $\vec{s}_{\alpha}$ be local oriented orthonormal frames over $U_{\alpha}$. On the overlap, we express $\vec{s}_{\alpha}=$ $g_{\alpha \beta} \vec{s}_{\beta}$ where $g_{\alpha \beta}: U_{\alpha} \cap U_{\beta} \rightarrow \mathrm{SO}(2 k)$. These satisfy the cocycle condition:

$$
g_{\alpha \beta} g_{\beta \gamma} g_{\gamma \alpha}=I \quad \text { and } \quad g_{\alpha \alpha}=I
$$

The principal bundle $P_{\text {SO }}$ of oriented orthonormal frames has transition functions $g_{\alpha \beta}$ acting on $\mathrm{SO}(2 k)$ from the left.

A spin structure on $V$ is a lifting of the transition functions to $\operatorname{SPIN}(2 k)$ preserving the cocyle condition. If the lifting is denoted by $g^{\prime}$, then we assume:

$$
\rho\left(g_{\alpha \beta}^{\prime}\right)=g_{\alpha \beta}, \quad g_{\alpha \beta}^{\prime} g_{\beta \gamma}^{\prime} g_{\gamma \alpha}^{\prime}=I, \quad \text { and } \quad g_{\alpha \alpha}^{\prime}=I
$$

This is equivalent to constructing a principal $\operatorname{SPIN}(2 k)$ bundle $P_{\text {SPIN }}$ together with a double covering map $\rho: P_{\text {SPIN }} \rightarrow P_{\text {SO }}$ which preserves the group action-i.e.,

$$
\rho\left(x \cdot g^{\prime}\right)=\rho(x) \cdot \rho\left(g^{\prime}\right)
$$

The transition functions of $P_{\text {SPIN }}$ are just the $g_{\alpha \beta}^{\prime}$ acting on the left.
Attempting to find a spin structure on $V$ is equivalent to taking a square root in a certain sense which we will make clear later. There is an obstruction to defining a spin structure which is similar to that to defining an orientation on $V$. These obstructions are $\mathbf{Z}_{2}$ characteristic classes called the Stieffel-Whitney classes and are most easily defined in terms of Čech cohomology. To understand these obstructions better, we review the construction briefly.

We fix a Riemannian structure on $T(M)$ to define a notion of distance. Geodesics on $M$ are curves which locally minimize distance. A set $U$ is said to be geodesically convex if (a) given $x, y \in U$ there exists a unique geodesic in $M$ joining $x$ to $y$ with $d(x, y)=\operatorname{length}(\gamma)$ and (b) if $\gamma$ is any such geodesic then $\gamma$ is actually contained in $U$. It is immediate that the intersection of geodesically convex sets is again geodesically convex and that every geodesically convex set is contractible.

It is a basic theorem of Riemannian geometry that there exist open covers of $M$ by geodesically convex sets. A cover $\left\{U_{\alpha}\right\}$ is said to be simple if the intersection of any number of sets of the cover is either empty or is contractible. A cover of $M$ by open geodesically convex sets is a simple cover.

We fix such a simple cover hence forth. Since $U_{\alpha}$ is contractible, any vector bundle over $M$ is trivial over $U_{\alpha}$. Let $\mathbf{Z}_{2}$ be the multiplicative group $\{ \pm 1\}$. A Čech $j$-cochain is a function $f\left(\alpha_{0}, \ldots, \alpha_{j}\right) \in \mathbf{Z}_{2}$ defined for $j+1$-tuples of indices where $U_{\alpha_{0}} \cap \cdots \cap U_{\alpha_{j}} \neq \emptyset$ which is totally symmetric-i.e.,

$$
f\left(\alpha_{\sigma(0)}, \ldots, \alpha_{\sigma(j)}\right)=f\left(\alpha_{0}, \ldots, \alpha_{j}\right)
$$

for any permutation $\sigma$. If $C^{j}\left(M ; \mathbf{Z}_{2}\right)$ denotes the multiplicative group of all such functions, the coboundary $\delta: C^{j}\left(M, \mathbf{Z}_{2}\right) \rightarrow C^{j+1}\left(M ; \mathbf{Z}_{2}\right)$ is defined by:

$$
(\delta f)\left(\alpha_{0}, \ldots, \alpha_{j+1}\right)=\prod_{i=0}^{j+1} f\left(\alpha_{0}, \ldots, \widehat{\alpha_{i}}, \ldots, \alpha_{j+1}\right)
$$

The multiplicative identity of $C^{j}\left(M ; \mathbf{Z}_{2}\right)$ is the function 1 and it is an easy combinatorial exercise that $\delta^{2} f=1$. For example, if $f_{0} \in C^{0}\left(M ; \mathbf{Z}_{2}\right)$ and if $f_{1} \in C^{1}\left(M ; \mathbf{Z}_{2}\right)$, then:

$$
\begin{aligned}
\delta\left(f_{0}\right)\left(\alpha_{0}, \alpha_{1}\right) & =f\left(\alpha_{1}\right) f\left(\alpha_{0}\right) \\
\delta\left(f_{1}\right)\left(\alpha_{0}, \alpha_{1}, \alpha_{2}\right) & =f\left(\alpha_{1}, \alpha_{2}\right) f\left(\alpha_{0}, \alpha_{1}\right) f\left(\alpha_{o}, \alpha_{2}\right)
\end{aligned}
$$

$\mathbf{Z}_{2}$ is a particularly simple coefficient group to work with since every element is its own inverse; in defining the Čech cohomology with coefficients
in other abelian groups, more care must be taken with the signs which arise.

Let $H^{j}\left(M ; \mathbf{Z}_{2}\right)=\mathrm{N}\left(\delta_{j}\right) / \mathrm{R}\left(\delta_{j-1}\right)$ be the cohomology group; it is an easy exercise to show these groups are independent of the particular simple cover chosen. There is a ring structure on $H^{*}\left(M ; \mathbf{Z}_{2}\right)$, but we shall only use the "additive" group structure (which we shall write multiplicatively).

Let $V$ be a real vector bundle, not necessarily orientable. Since the $U_{\alpha}$ are contractible, $V$ is trivial over the $U$ and we can find a local orthonormal frame $\vec{s}_{\alpha}$ for $V$ over $U_{\alpha}$. We let $\vec{s}_{\alpha}=g_{\alpha \beta} \vec{s}_{\beta}$ and define the 1-cochain:

$$
f(\alpha \beta)=\operatorname{det}\left(g_{\alpha \beta}\right)= \pm 1
$$

This is well defined since $U_{\alpha} \cap U_{\beta}$ is contractible and hence connected. Since $f(\alpha, \beta)=f(\beta, \alpha)$, this defines an element of $C^{1}\left(M ; \mathbf{Z}_{2}\right)$. Since the $\left\{g_{\alpha \beta}\right\}$ satisfy the cocycle condition, we compute:

$$
\delta f(\alpha, \beta, \gamma)=\operatorname{det}\left(g_{\alpha \beta} g_{\beta \gamma} g_{\gamma \alpha}\right)=\operatorname{det}(I)=1
$$

so $\delta(f)=1$ and $f$ defines an element of $H^{1}\left(M ; \mathbf{Z}_{2}\right)$. If we replace $\vec{s}_{\alpha}$ by $\vec{s}_{\alpha}^{\prime}=h_{\alpha} \vec{s}_{\alpha}$ the new transition functions become $g_{\alpha \beta}^{\prime}=h_{\alpha} g_{\alpha \beta} h_{\beta}^{-1}$ so if $f_{0}(\alpha)=\operatorname{det}\left(h_{\alpha}\right)$,

$$
f^{\prime}(\alpha, \beta)=\operatorname{det}\left(h_{\alpha} g_{\alpha \beta} h_{\beta}^{-1}\right)=\operatorname{det}\left(h_{\alpha}\right) f(\alpha, \beta) \operatorname{det}\left(h_{\beta}\right)=\delta\left(f_{0}\right) f
$$

and $f$ changes by a coboundary. This proves the element in cohomology defined by $f$ is independent of the particular frame chosen and we shall denote this element by $w_{1}(V) \in H^{1}\left(M ; \mathbf{Z}_{2}\right)$.

If $V$ is orientable, we can choose frames so $\operatorname{det}\left(g_{\alpha \beta}\right)=1$ and thus $w_{1}(V)=1$ represents the trivial element in cohomology. Conversely, if $w_{1}(V)$ is trivial, then $f=\delta f_{0}$. If we choose $h_{\alpha}$ so $\operatorname{det}\left(h_{\alpha}\right)=f_{0}(\alpha)$, then the new frames $\vec{s}_{\alpha}^{\prime}=h_{\alpha} \vec{s}_{\alpha}$ will have transition functions with $\operatorname{det}\left(g_{\alpha \beta}^{\prime}\right)=1$ and define an orientation of $V$. Thus $V$ is orientable if and only if $w_{1}(V)$ is trivial and $w_{1}(V)$, which is called the first Stieffel-Whitney class, measures the obstruction to orientability.

If $V$ is orientable, we restrict henceforth to oriented frames. Let $\operatorname{dim} V$ $=2 k$ be even and let $g_{\alpha \beta} \in \mathrm{SO}(2 k)$ be the transition functions. We choose any lifting $\tilde{g}_{\alpha \beta}$ to $\operatorname{SPIN}(2 k)$ so that:

$$
\rho\left(\tilde{g}_{\alpha \beta}\right)=g_{\alpha \beta} \quad \text { and } \quad \tilde{g}_{\alpha \beta} \tilde{g}_{\beta \alpha}=I
$$

since the $U_{\alpha}$ are contractible such lifts always exist. We have $g_{\alpha \beta} g_{\beta \gamma} g_{\gamma \alpha}$ $=I$ so $\rho\left(\tilde{g}_{\alpha \beta} \tilde{g}_{\beta \gamma} \tilde{g}_{\gamma \alpha}\right)=I$ and hence $\tilde{g}_{\alpha \beta} \tilde{g}_{\beta \gamma} \tilde{g}_{\gamma \alpha}= \pm I=f(\alpha, \beta, \gamma) I$ where $f(\alpha, \beta, \gamma) \in \mathbf{Z}_{2} . V$ admits a spin structure if and only if we can choose the lifting so $f(\alpha, \beta, \gamma)=1$. It is an easy combinatorial exercise to show
that $f$ is symmetric and that $\delta f=0$. Furthermore, if we change the choice of the $\vec{s}_{\alpha}$ or change the choice of lifts, then $f$ changes by a coboundary. This implies $f$ defines an element $w_{2}(V) \in H^{2}\left(M ; \mathbf{Z}_{2}\right)$ independent of the choices made. $w_{2}$ is called the second Stieffel-Whitney class and is trivial if and only if $V$ admits a spin structure.

We suppose $w_{1}(V)=w_{2}(V)=1$ are trivial so $V$ is orientable and admits a spin structure. Just as there are two orientations which can be chosen on $V$, there can be several possible inequivalent spin structures (i.e., several non-isomorphic principal bundles $\left.P_{\text {SPIN }}\right)$. It is not difficult to see that inequivalent spin structures are parametrized by representations of the fundamental group $\pi_{1}(M) \rightarrow \mathbf{Z}_{2}$ just as inequivalent orientations are parametrized by maps of the components of $M$ into $\mathbf{Z}_{2}$.

To illustrate the existence and non-existence of spin structures on bundles, we take $M=S^{2}$. $S^{2}=\mathbf{C} P_{1}$ is the Riemann sphere. There is a natural projection from $\mathbf{C}^{2}-0$ to $S^{2}$ given by sending $(x, w) \mapsto z / w ; S^{2}$ is obtained by identifying $(z, w)=(\lambda z, \lambda w)$ for $(z, w) \neq(0,0)$ and $\lambda \neq 0$. There are two natural charts for $S^{2}$ :

$$
U_{1}=\{z:|z| \leq 1\} \quad \text { and } \quad U_{2}=\{z:|z| \geq 1\} \cup\{\infty\}
$$

where $w=1 / z$ gives coordinates on $U_{2}$.
$\mathbf{C} P_{1}$ is the set of lines in $\mathbf{C}^{2}$; we let $L$ be the natural line bundle over $S^{2}$; this is also refered to as the tautological line bundle. We have natural sections to $L$ over $U_{i}$ defined by:

$$
s_{1}=(z, 1) \text { over } U_{1} \quad \text { and } \quad s_{2}=(1, w) \text { over } U_{2}
$$

these are related by the transition function:

$$
s_{1}=z s_{2}
$$

so on $U_{1} \cap U_{2}$ we have $g_{12}=e^{i \theta}$. The double cover $\operatorname{SPIN}(2) \rightarrow \mathrm{SO}(2)$ is defined by $\theta \mapsto 2 \theta$ so this bundle does not have a spin structure since this transition function cannot be lifted to SPIN(2).

This cover is not a simple cover of $S^{2}$ so we construct the cover given in the following
$V_{4}$

(where we should "fatten up" the picture to have an open cover). If we have sections $\vec{s}_{i}$ to $L$ over $V_{i}$, then we can choose the transition functions so $\vec{s}_{j}=e^{i \theta} \vec{s}_{4}$ for $1 \leq j \leq 3$ and $\vec{s}_{1}=\vec{s}_{2}=\vec{s}_{3}$. We let

$$
\begin{array}{rll}
0 \leq \theta \leq 2 \pi / 3 & \text { parametrize } & V_{1} \cap V_{4} \\
2 \pi / 3 \leq \theta \leq 4 \pi / 3 & \text { parametrize } & V_{2} \cap V_{4} \\
4 \pi / 3 \leq \theta \leq 2 \pi & \text { parametrize } & V_{3} \cap V_{4}
\end{array}
$$

and define:

$$
\tilde{g}_{14}=e^{i \theta / 2}, \quad \tilde{g}_{24}=e^{i \theta / 2}, \quad \tilde{g}_{34}=e^{i \theta / 2}
$$

and all the $\tilde{g}_{j k}=1$ for $1 \leq j, k \leq 3$ then we compute:

$$
\tilde{g}_{13} \tilde{g}_{34} \tilde{g}_{41}=e^{0} e^{i \pi}=-1
$$

so the cochain defining $w_{2}$ satisfies:

$$
f(1,3,4)=-1 \quad \text { with } \quad f(i, j, k)=1 \quad \text { if } i \text { or } j \text { or } k=2 .
$$

We compute this is non-trivial in cohomology as follows: suppose $f=\delta f_{1}$. Then:

$$
\begin{aligned}
& 1=f(1,2,3)=f_{1}(1,2) f_{1}(1,3) f_{1}(2,3) \\
& 1=f(1,2,4)=f_{1}(1,2) f_{1}(1,4) f_{1}(2,4) \\
& 1=f(2,3,4)=f_{1}(2,3) f_{1}(2,4) f_{1}(3,4)
\end{aligned}
$$

and multiplying together:

$$
\begin{aligned}
1 & =f_{1}(1,2)^{2} f_{1}(2,3)^{2} f_{1}(2,4)^{2} f_{1}(1,4) f_{1}(1,3) f_{1}(3,4) \\
& =f_{1}(1,4) f_{1}(1,3) f_{1}(3,4)=f(1,3,4)=-1
\end{aligned}
$$

which is a contradiction. Thus we have computed combinatorially that $w_{2} \neq 1$ is non-trivial.

Next we let $V=T(M)$ be the real tangent bundle. We identify $\mathrm{U}(1)$ with $\mathrm{SO}(2)$ to identify $T(M)$ with $T_{C}(M)$. Since $w=1 / z$ we have:

$$
\frac{d}{d z}=\frac{d w}{d z} \frac{d}{d w}=-z^{-2} \frac{d}{d w}
$$

so the transition function on the overlap is $-e^{-2 i \theta}$. The minus sign can be eliminated by using $-\frac{d}{d w}$ instead of $\frac{d}{d w}$ as a section over $U_{2}$ to make the transition function be $e^{-2 i \theta}$. Since the double cover of $\operatorname{SPIN}(2) \rightarrow \mathrm{SO}(2)$ is given by $\theta \mapsto 2 \theta$, this transition function lifts and $T(M)$ has a spin structure. Since $S^{2}$ is simply connected, the spin structure is unique. This
also proves $T_{C}(M)=L^{*} \otimes L^{*}$ since the two bundles have the same transition functions.

There is a natural inclusion of $\operatorname{SPIN}(V)$ and $\operatorname{SPIN}(W)$ into subgroups of $\operatorname{SPIN}(V \oplus W)$ which commute. This induces a map from $\operatorname{SPIN}(V) \times$ $\operatorname{SPIN}(W) \rightarrow \operatorname{SPIN}(V \oplus W)$. Using this map, it is easy to compute $w_{2}(V \oplus$ $W)=w_{2}(V) w_{2}(W)$. To define $w_{2}(V)$, we needed to choose a fixed orientation of $V$. It is not difficult to show that $w_{2}(V)$ is independent of the orientation chosen and of the fiber metric on $V$. (Our definition does depend upon the fact that $V$ is orientable. Although it is possible to define $w_{2}(V)$ more generally, the identity $w_{2}(V \oplus W)=w_{2}(V) w_{2}(W)$ fails if $V$ and $W$ are not orientable in general). We summarize the properties of $w_{2}$ we have derived:

Lemma 3.3.1. Let $V$ be a real oriented vector bundle and let $w_{2}(V) \in$ $H^{2}\left(M ; \mathbf{Z}_{2}\right)$ be the second Stieffel-Whitney class. Then:
(a) $w_{2}(V)=1$ represents the trivial cohomology class if and only if $V$ admits a spin structure.
(b) $w_{2}(V \oplus W)=w_{2}(V) w_{2}(W)$.
(c) If $L$ is the tautological bundle over $S^{2}$ then $w_{2}(L)$ is non-trivial.

We emphasize that (b) is written in multiplicative notation since we have chosen the multiplicative version of $\mathbf{Z}_{2} . \quad w_{2}$ is also functorial under pull-backs, but as we have not defined the Cech cohomology as a functor, we shall not discuss this property. We also note that the Stieffel-Whitney classes can be defined in general; $w(V)=1+w_{1}(V)+\cdots \in H^{*}\left(M ; \mathbf{Z}_{2}\right)$ is defined for any real vector bundle and has many of the same properties that the Chern class has for complex bundles.

We can use Lemma 3.3.1 to obtain some other results on the existence of spin structures:

Lemma 3.3.2.
(a) If $M=S^{m}$ then any bundle over $S^{m}$ admits a spin structure for $m \neq 2$.
(b) If $M=\mathbf{C} P_{k}$ and if $L$ is the tautological bundle over $M$, then $L$ does not admit a spin structure.
(c) If $M=\mathbf{C} P_{k}$ and if $V=T(M)$ is the tangent space, then $V$ admits a spin structure if and only if $k$ is odd.
(d) If $M=\mathbf{Q} P_{k}$ is quaternionic projective space, then any bundle over $M$ admits a spin structure.
Proof: $H^{2}\left(M ; \mathbf{Z}_{2}\right)=\{1\}$ in (a) and (d) so $w_{2}$ must represent the trivial element. To prove (b), we suppose $L$ admits a spin structure. $S^{2}$ is embedded in $\mathbf{C} P_{k}$ for $k \geq 2$ and the restriction of $L$ to $S^{2}$ is the tautological bundle over $S^{2}$. This would imply $L$ admits a spin structure over $S^{2}$ which is false. Finally, we use the representation:

$$
T_{c}\left(\mathbf{C} P_{j}\right) \oplus 1=L^{*} \oplus \cdots \oplus L^{*}
$$

so that

$$
T\left(\mathbf{C} P_{j}\right) \oplus 1^{2}=\underbrace{\left(L_{r}^{*}\right) \oplus \cdots \oplus\left(L_{r}^{*}\right)}_{j+1 \text { times }}
$$

where $L_{r}^{*}$ denotes the real vector bundle obtained from $L^{*}$ by forgetting the complex structure. Then $w_{2}\left(L_{r}^{*}\right)=w_{2}\left(L_{r}\right)$ since these bundles are isomorphic as real bundles. Furthermore:

$$
w_{2}\left(T\left(\mathbf{C} P_{j}\right)\right)=w_{2}\left(L_{r}\right)^{j+1}
$$

and this is the trivial class if $j+1$ is even (i.e., $j$ is odd) while it is $w_{2}\left(L_{r}\right)$ and non-trivial if $j+1$ is odd (i.e., $j$ is even).

Let $M$ be a manifold and let $V$ be a real vector bundle over $M$. If $V$ admits a spin structure, we shall say $V$ is spin and not worry for the moment about the lack of uniqueness of the spin structure. This will not become important until we discuss Lefschetz fixed point formulas.

If $V$ is spin, we define the bundles $\Delta^{ \pm}(V)$ to have transition functions $\Delta^{ \pm}\left(\tilde{g}_{\alpha \beta}\right)$ where we apply the representation $\Delta^{ \pm}$to the lifted transition functions. Alternatively, if $P_{\text {SPIN }}$ is the principal spin bundle defining the spin structure, we define:

$$
\Delta^{ \pm}(V)=P_{\text {SPIN }} \otimes_{\Delta^{ \pm}} \Delta^{ \pm}
$$

where the tensor product across a representation is defined to be $P_{\text {SPIN }} \times$ $\Delta^{ \pm}$module the identification $(p \cdot g) \times z=p \times\left(\Delta^{ \pm}(g) z\right)$ for $p \in P_{\text {SPIN }}$, $g \in P_{\text {SPIN }}(2 k), z \in \Delta^{ \pm}$. (The slight confusion of notation is caused by our convention of using the same symbol for the representation and the representation space.)

Let $b_{ \pm}$be a fixed unitary frame for $\Delta^{ \pm}$. If $\vec{s}$ is a local oriented orthonormal frame for $V$, we let $\tilde{s}_{i}$ be the two lifts of $\vec{s} \in P_{\text {SO }}$ to $P_{\text {SPIN }} . \tilde{s}_{1}=-\tilde{s}_{2}$ but there is no natural way to distinguish these lifts, although the pair is cannonically defined. If $\nabla$ is a Riemannian connection on $V$, let $\nabla_{\vec{s}}=\omega \vec{s}$ be the connection 1-form. $\omega$ is a skew symmetric matrix of 1 -forms and is a 1 -form valued element of the Lie algebra of $\mathrm{SO}(2 k)$. Since the Lie algebra of $\mathrm{SO}(2 k)$ and $\operatorname{SPIN}(2 k)$ coincide, we can also regard $\omega$ as an element which is 1 -form valued of the Lie algebra of $\operatorname{SPIN}(2 k)$ and let $\Delta^{ \pm}(\omega)$ act on $\Delta^{ \pm}(V)$. We define bases $\tilde{s}_{i} \otimes b_{ \pm}$for $\Delta^{ \pm}(V)$ and define:

$$
\nabla\left(\tilde{s}_{i} \otimes b_{ \pm}\right)=\tilde{s}_{i} \otimes \Delta^{ \pm}(\omega) b_{ \pm}
$$

to define a natural connection on $\Delta^{ \pm}(V)$. Since the same connection is defined whether $\tilde{s}_{1}$ or $\tilde{s}_{2}=-\tilde{s}_{1}$ is chosen, the $\mathbf{Z}_{2}$ ambiguity is irrelevant and $\nabla$ is well defined.

Lemma 3.2.4 and 3.2.5 extend to:

Lemma 3.3.3. Let $V$ be a real oriented Riemannian vector bundle of even fiber dimension $v$. Suppose $V$ admits a spin structure and let $\Delta^{ \pm}(V)$ be the half-spin bundles. Then:
(a) $\Lambda(V) \simeq \Delta(V) \otimes \Delta(V)$.
(b) $\left(\Lambda^{+}-\Lambda^{-}\right)(V) \simeq\left(\Delta^{+}-\Delta^{-}\right) \otimes\left(\Delta^{+}-\Delta^{-}\right)(V)$.
(c) $\left(\Lambda^{\mathrm{e}}-\Lambda^{\mathrm{o}}\right)(V) \simeq(-1)^{v / 2}\left(\Delta^{+}-\Delta^{-}\right) \otimes\left(\Delta^{+}-\Delta^{-}\right)(V)$.
(d) If $v=2$ and if $V$ is the underlying real bundle of a complex bundle $V_{c}$ then $V_{c} \simeq \Delta^{-} \otimes \Delta^{-}$and $V_{c}^{*} \simeq \Delta^{+} \otimes \Delta^{+}$.
(e) If $\nabla$ is a Riemannian connection on $V$ and if we extend $\nabla$ to $\Delta^{ \pm}(V)$, then the isomorphisms given are unitary isomorphisms which preserve $\nabla$.

Spinors are multiplicative with respect to products. It is immediate from the definitions we have given that:

Lemma 3.3.4. Let $V_{i}$ be real oriented Riemannian vector bundles of even fiber dimensions $v_{i}$ with given spin structures. Let $V_{1} \oplus V_{2}$ have the natural orientation and spin structure. Then:

$$
\left(\Delta^{+}-\Delta^{-}\right)\left(V_{1} \oplus V_{2}\right) \simeq\left\{\left(\Delta^{+}-\Delta^{-}\right)\left(V_{1}\right)\right\} \otimes\left\{\left(\Delta^{+}-\Delta^{-}\right)\left(V_{2}\right)\right\}
$$

If $\nabla_{i}$ are Riemannian connections on $V_{i}$ and if we define the natural induced connections on these bundles, then the isomorphism is unitary and preserves the connections.

A spin structure always exists locally since the obstruction to a spin structure is global. Given any real oriented Riemannian bundle $V$ of even fiber dimension with a fixed Riemannian connection $\nabla$, we define $\Lambda(V)$ and $\Delta(V)=\Delta^{+}(V) \oplus \Delta^{-}(V)$. With the natural metrics and connections, we relate the connection 1-forms and curvatures of these 3 bundles:
Lemma 3.3.5. Let $\left\{e_{i}\right\}$ be a local oriented orthonormal frame for $V$. We let $\nabla e_{j}=\omega_{j k} s_{k}$ represent the connection 1-form and $\Omega e_{j}=\Omega_{j k} e_{k}$ be the curvature matrix for $\Omega_{j k}=d \omega_{j k}-\omega_{j l} \wedge \omega_{l k}$. Let $\Lambda e$ and $\Delta e$ denote the natural orthonormal frames on $\Lambda(V)$ and $\Delta(V)$. Relative to these frames we compute the connection 1-forms and curvature matrices of $\Lambda(V)$ and $\Delta(V)$ by:

$$
\begin{array}{ll}
\omega_{\Lambda}=\omega_{j k} \operatorname{ext}\left(e_{k}\right) \operatorname{int}\left(e_{j}\right) & \\
\Omega_{\Lambda}=\Omega_{j k} \operatorname{ext}\left(e_{k}\right) \operatorname{int}\left(e_{j}\right)  \tag{b}\\
\omega_{\Delta}=\frac{1}{4} \omega_{j k} e_{j} * e_{k} & \\
\Omega_{\Delta}=\frac{1}{4} \Omega_{j k} e_{j} * e_{k}
\end{array}
$$

Proof: We sum over repeated indices in these expressions. We note (a) is true by definition on $\Lambda^{1}(V)=V$. Since both $\omega$ and $\Omega$ extend to act as derivations on the exterior algebra, this implies (a) is true on forms of all degree.

Let $\operatorname{so}(n)=\left\{A \in n \times n\right.$ real matrices with $\left.A+A^{t}=0\right\}$ be the Lie algebra of $\operatorname{SO}(n)$. This is also the Lie algebra of $\operatorname{SPIN}(n)$ so we must identify this
with an appropriate subset of the Clifford algebra in order to prove (b). We choose a representative element of $\operatorname{so}(n)$ and define

$$
A e_{1}=e_{2}, \quad A e_{2}=-e_{1}, \quad A e_{j}=0 \quad \text { for } j>2
$$

i.e.,

$$
A_{12}=1, \quad A_{21}=-1, \quad A_{j k}=0 \quad \text { otherwise }
$$

If we let $g(t) \in \mathrm{SO}(n)$ be defined by

$$
\begin{aligned}
g(t) e_{1} & =(\cos t) e_{1}+(\sin t) e_{2}, \\
g(t) e_{2} & =(\cos t) e_{2}-(\sin t) e_{1}, \\
g(t) e_{j} & =e_{j} \quad j>2,
\end{aligned}
$$

then $g(0)=I$ and $g^{\prime}(0)=A$. We must lift $g(t)$ from $\operatorname{SO}(n)$ to $\operatorname{SPIN}(n)$. Define:

$$
\begin{aligned}
h(t) & =\left(\cos (t / 4) e_{1}+\sin (t / 4) e_{2}\right)\left(-\cos (t / 4) e_{1}+\sin (t / 4) e_{2}\right) \\
& =\cos (t / 2)+\sin (t / 2) e_{1} e_{2} \in \operatorname{SPIN}(n)
\end{aligned}
$$

Then $\rho(h) \in \operatorname{SO}(n)$ is defined by:

$$
\rho(h) e_{j}=\left(\cos (t / 2)+\sin (t / 2) e_{1} e_{2}\right) e_{j}\left(\cos (t / 2)-\sin (t / 2) e_{1} e_{2}\right)
$$

so that $\rho(h) e_{j}=e_{j}$ for $j>2$. We compute:

$$
\begin{aligned}
\rho(h) e_{1} & =\left(\cos (t / 2) e_{1}+\sin (t / 2) e_{2}\right)\left(\cos (t / 2)-\sin (t / 2) e_{1} e_{2}\right) \\
& =\left(\cos ^{2}(t / 2)-\sin ^{2}(t / 2)\right) e_{1}+2 \sin (t / 2) \cos (t / 2) e_{2} \\
\rho(h) e_{2} & =\left(\cos (t / 2) e_{2}-\sin (t / 2) e_{1}\right)\left(\cos (t / 2)-\sin (t / 2) e_{1} e_{2}\right) \\
& =\left(\cos ^{2}(t / 2)-\sin ^{2}(t / 2)\right) e_{2}-2 \sin (t / 2) \cos (t / 2) e_{1}
\end{aligned}
$$

so that $\rho(h)=g$. This gives the desired lift from $\operatorname{SO}(n)$ to $\operatorname{SPIN}(n)$. We differentiate to get

$$
h^{\prime}(0)=\frac{1}{2} e_{1} e_{2}=\frac{1}{4} A_{j k} e_{j} * e_{k} .
$$

This gives the lift of a matrix in this particular form. Since the whole Lie algebra is generated by elements of this form, it proves that the lift of $A_{j k}$ in general is given by $\frac{1}{4} A_{j k} e_{j} * e_{k}$. Since $\operatorname{SPIN}(n)$ acts on the Clifford algebra by Clifford multiplication on the left, this gives the action of the curvature and connection 1 -form on the spin representations and completes the proof.

We can define $\operatorname{ch}\left(\Delta^{ \pm}(V)\right)$ as SO characteristic forms. They can be expressed in terms of the Pontrjagin and Euler forms. We could use the explicit representation given by Lemma 3.3.5 to compute these forms; it is most easy, however, to compute in terms of generating functions. We introduce formal variables $\left\{x_{j}\right\}$ for $1 \leq j \leq(\operatorname{dim} V) / 2$ so that

$$
p(V)=\prod_{j}\left(1+x_{j}^{2}\right) \quad \text { and } \quad e(V)=\prod_{j} x_{j}
$$

All these computations are really induced from the corresponding matrix identities and the $\left\{x_{j}\right\}$ arise from putting a skew-symmetric matrix $A$ in block form.

Lemma 3.3.6. Let $V$ be an oriented real Riemannian vector bundle of dimension $n \equiv 0$ (2). Let $\operatorname{ch}\left(\Delta^{ \pm}(V)\right)$ be a real characteristic form; these are well defined even if $V$ does not admit a global spin structure. Then:
(a) $\operatorname{ch}\left(\Delta^{+}(V)\right)+\operatorname{ch}\left(\Delta^{-}(V)\right)=\operatorname{ch}(\Delta(V))=\prod_{j}\left\{e^{x_{j} / 2}+e^{-x_{j} / 2}\right\}$.
(b)

$$
\begin{aligned}
(-1)^{n / 2}\left\{\operatorname{ch}\left(\Delta^{+}(V)\right)-\operatorname{ch}\left(\Delta^{-}(V)\right)\right\} & =\prod_{j}\left\{e^{x_{j} / 2}-e^{-x_{j} / 2}\right\} \\
& =e(V)(1+\text { higher order terms })
\end{aligned}
$$

Proof: This is an identity among invariant polynomials in the Lie algebra $\operatorname{so}(n)$. We may therefore restrict to elements of the Lie algebra which split in block diagonal form. Using the multiplicative properties of the Chern character and Lemma 3.3.4, it suffices to prove this lemma for the special case that $n=2$ so $\Delta^{ \pm}(V)$ are complex line bundles. Let $V_{c}$ be a complex line bundle and let $V$ be the underlying real bundle. Then $x_{1}=x=c_{1}\left(V_{c}\right)=e(V)$. Since:

$$
V_{c}=\Delta^{-} \otimes \Delta^{-} \quad \text { and } \quad V_{c}^{*}=\Delta^{+} \otimes \Delta^{+}
$$

we conclude:

$$
x=2 c_{1}\left(\Delta^{-}\right) \quad \text { and } \quad-x=2 c_{1}\left(\Delta^{+}\right)
$$

which shows (a) and the first part of (b). We expand:

$$
e^{x / 2}-e^{-x / 2}=x+\frac{1}{24} x^{3}+\cdots
$$

to see $\operatorname{ch}\left(\Delta^{-}\right)-\operatorname{ch}\left(\Delta^{+}\right)=e(V)\left(1+\frac{1}{24} p_{1}(V)+\cdots\right)$ to complete the proof of (b).

There is one final calculation in characteristic classes which will prove helpful. We defined $L$ and $\hat{A}$ by the generating functions:

$$
\begin{aligned}
L(x) & =\prod_{j} \frac{x_{j}}{\tanh x_{j}} \\
& =\prod_{j} x_{j} \frac{e^{x_{j}}+e^{-x_{j}}}{e^{x_{j}}-e^{-x_{j}}} \\
\hat{A}(x) & =\prod_{j} \frac{x_{j}}{\sinh \left(x_{j} / 2\right)}=\prod_{j} \frac{x_{j}}{e^{x_{j} / 2}-e^{-x_{j} / 2}}
\end{aligned}
$$

Lemma 3.3.7. Let $V$ be an oriented Riemannian real vector bundle of dimension $n=2 n_{0}$. If we compute the component of the differential form which is in $\Lambda^{n}\left(T^{*} M\right)$ then:

$$
\{L(V))\}_{n}=\{\operatorname{ch}(\Delta(V)) \wedge \hat{A}(V)\}_{n}
$$

Proof: It suffices to compute the component of the corresponding symmetric functions which are homogeneous of order $n_{0}$. If we replace $x_{j}$ by $x_{j} / 2$ then:

$$
\begin{aligned}
\left\{L\left(x_{j}\right)\right\} n_{0} & =\left\{2^{n_{0}} L\left(x_{j} / 2\right)\right\}_{n_{0}} \\
& =\left\{\prod_{j} x_{j} \tanh \left(x_{j} / 2\right)\right\}_{n_{0}} \\
& =\left\{\prod_{j} x_{j} \frac{e^{x / 2}+e^{-x / 2}}{e^{x_{j} / 2}-e^{-x_{j} / 2}}\right\}_{n_{0}} \\
& =\{\operatorname{ch}(\Delta(V)) \wedge \hat{A}(V)\}_{n_{0}}
\end{aligned}
$$

which completes the proof.

### 3.4. The Spin Complex.

We shall use the spin complex chiefly as a formal construction to link the de Rham, signature, and Dolbeault complexes. Let $M$ be a Riemannian manifold of even dimension $m$. Let $T(M)$ be the real tangent space. We assume that $M$ is orientable and that $T(M)$ admits a spin structure. We let $\Delta^{ \pm}(M)$ be the half-spin representations. There is a natural map given by Lemma 3.2.3 from the representations of

$$
T^{*}(M) \rightarrow \operatorname{HOM}\left(\Delta^{ \pm}(M), \Delta^{\mp}(M)\right)
$$

which we will call $c(\xi)$ (since it is essentially Clifford multiplication) such that $c(\xi)^{2}=-|\xi|^{2} I$. We extend the Levi-Civita connection to act naturally on these bundles and define the spin complex by the diagram:

$$
A^{ \pm}: C^{\infty}\left(\Delta^{ \pm}(M)\right) \xrightarrow{\nabla} C^{\infty}\left(T^{*} M \otimes \Delta^{ \pm}(M)\right) \xrightarrow{c} C^{\infty}\left(\Delta^{\mp}(M)\right)
$$

to be the operator with leading symbol $c .\left(A^{+}\right)^{*}=A^{-}$and $A^{ \pm}$is elliptic since $c(\xi)^{2}=-|\xi|^{2} I$; this operator is called the Dirac operator.

Let $V$ be a complex bundle with a Riemannian connection $\nabla$. We define the spin complex with coefficients in $\nabla$ by using the diagram:

$$
\begin{aligned}
A_{V}^{ \pm}: C^{\infty}\left(\Delta^{ \pm}(M) \otimes V\right) \rightarrow C^{\infty}\left(T^{*} M \otimes \Delta^{ \pm}(M) \otimes V\right) & \\
& \xrightarrow{c \otimes 1} C^{\infty}\left(\Delta^{\mp}(M) \otimes V\right) .
\end{aligned}
$$

This is completely analogous to the signature complex with coefficients in $V$. We define:

$$
\operatorname{index}(V, \operatorname{spin})=\operatorname{index}\left(A_{V}^{+}\right)
$$

a priori this depends on the particular spin structure chosen on $V$, but we shall show shortly that it does not depend on the particular spin structure, although it does depend on the orientation. We let

$$
a_{n}^{\text {spin }}(x, V)
$$

denote the invariants of the heat equation. If we reverse the orientation of $M$, we interchange the roles of $\Delta^{+}$and of $\Delta^{-}$so $a_{n}^{\text {spin }}$ changes sign. This implies $a_{n}^{\text {spin }}$ can be regarded as an invariantly defined $m$-form since the scalar invariant changes sign if we reverse the orientation.

Let $X$ be an oriented coordinate system. We apply the Gramm-Schmidt process to the coordinate frame to construct a functorial orthonormal frame $\vec{s}(X)$ for $T(M)=T^{*}(M)$. We lift this to define two local sections $\vec{s}_{i}(X)$ to the principal bundle $P_{\text {SPIN }}$ with $\vec{s}_{1}(X)=-\vec{s}_{2}(X)$. There is, of course, no cannonical way to prefer one over the other, but the pair is invariantly defined. Let $b^{ \pm}$be fixed bases for the representation space $\Delta^{ \pm}$and let
$b_{i}^{ \pm}(X)=\vec{s}_{i}(X) \otimes b^{ \pm}$provide local frames for $\Delta^{ \pm}(M)$. If we fix $i=1$ or $i=2$, the symbol of the spin complex with coefficients in $V$ can be functorially expressed in terms of the 1-jets of the metric on $M$ and in terms of the connection 1 -form on $V$. The leading symbol is given by Clifford multiplication; the $0^{\text {th }}$ order term is linear in the 1-jets of the metric on $M$ and in the connection 1-form on $V$ with coefficients which are smooth functions of the metric on $M$. If we replace $b_{1}^{ \pm}$by $b_{2}^{ \pm}=-b_{1}^{ \pm}$then the local representation of the symbol is unchanged since multiplication by -1 commutes with differential operators. Thus we may regard $a_{n}^{\text {spin }}(x, V) \in$ $\mathcal{R}_{m, n, m, \operatorname{dim} V}$ as an invariantly defined polynomial which is homogeneous of order $n$ in the jets of the metric and of the connection form on $V$ which is $m$-form valued.

This interpretation defines $a_{n}^{\text {spin }}(x, V)$ even if the base manifold $M$ is not spin. We can always define the spin complex locally as the $\mathbf{Z}_{2}$ indeterminacy in the choice of a spin structure will not affect the symbol of the operator. Of course, $\int_{M} a_{m}^{\text {spin }}(x, V)$ can only be given the interpretation of index ( $V$, spin) if $M$ admits a spin structure. In particular, if this integral is not an integer, $M$ cannot admit a spin structure.
$a_{n}^{\text {spin }}(x, V)$ is a local invariant which is not affected by the particular global spin structure chosen. Thus index ( $V$, spin) is independent of the particular spin structure chosen.

We use exactly the same arguments based on the theorem of the second chapter and the multiplicative nature of the twisted spin complex as were used to prove the Hirzebruch signature theorem to establish:

Lemma 3.4.1.
(a) $a_{n}^{\text {spin }}=0$ if $n<m$ and $a_{m}^{\text {spin }}$ is a characteristic form of $T(M)$ and of $V$.
(b) $\int_{M} a_{m}^{\text {spin }}(x, V)=\operatorname{index}(V$, spin $)$.
(c) There exists a characteristic form $\hat{A}^{\prime}$ of $T(M)$ in the form:

$$
\hat{A}^{\prime}=1+\hat{A}_{1}+\cdots
$$

which does not depend on the dimension of $M$ together with a universal constant $c$ such that

$$
a_{m}^{\mathrm{spin}}=\sum_{4 s+2 t=m} \hat{A}_{s}^{\prime} \wedge c^{t} c h(V)_{t} .
$$

We use the notation $\hat{A}^{\prime}$ since we have not yet shown it is the $A$-roof genus defined earlier. In proving the formula splits into this form, we do not rely on the uniqueness property of the second chapter, but rather on the multiplicative properties of the invariants of the heat equation discussed in the first chapter.

The spin complex has an intimate relation with both the de Rham and signature complexes:

Lemma 3.4.2. Let $U$ be an open contractible subset of $M$. Over $U$, we define the signature, de Rham, and spin complexes. Then:
(a) $\left(\Lambda^{\mathrm{e}}-\Lambda^{\mathrm{o}}\right) \otimes V \simeq(-1)^{m / 2}\left(\Delta^{+}-\Delta^{-}\right) \otimes\left(\Delta^{+}-\Delta^{-}\right) \otimes V$.
(b) $\left(\Lambda^{+}-\Lambda^{-}\right) \otimes V \simeq\left(\Delta^{+}-\Delta^{-}\right) \otimes\left(\Delta^{+}+\Delta^{-}\right) \otimes V$.
(c) These two isomorphisms preserve the unitary structures and the connections. They also commute with Clifford multiplication on the left.
(d) The natural operators on these complexes agree under this isomorphism.
Proof: (a)-(c) follow from previous results. The natural operators on these complexes have the same leading symbol and therefore must be the same since they are natural first order operators. This proves (d).

We apply this lemma in dimension $m=2$ with $V=$ the trivial bundle and $M=S^{2}$ :

$$
\begin{aligned}
\chi\left(S^{2}\right)=2 & =\operatorname{index}(d+\delta)=\operatorname{index}\left(\Delta^{-}-\Delta^{+}, \text {spin }\right) \\
& =\int_{S^{2}} c\left\{c_{1}\left(\Delta^{-}\right)-c_{1}\left(\Delta^{+}\right)\right\} \\
& =c \int_{S^{2}} e(T M)=2 c
\end{aligned}
$$

by Lemmas 2.3.1, 3.3.5, and 3.3.4. This establishes that the normalizing constant of Lemma 3.4.1 must be 1 so that:

$$
a_{m}^{\mathrm{spin}}=\sum_{4 s+2 t=m} \hat{A}^{\prime} \wedge c h_{t}(V)
$$

We apply this lemma to the twisted signature complex in dimension $m=2$ with $V$ a non-trivial line bundle over $S^{2}$ to conclude:

$$
\begin{aligned}
\operatorname{signature}\left(S^{2}, V\right) & =\operatorname{index}\left(\left(\Delta^{+} \oplus \Delta^{-}\right) \otimes V, \text { spin }\right) \\
& =\int_{S^{2}} \operatorname{ch}\left(\left(\Delta^{+} \oplus \Delta^{-}\right) \otimes V\right)=\int_{S^{2}} 2 \cdot c_{1}(V)
\end{aligned}
$$

This shows that the normalizing constant of Lemma 3.1.4 for the twisted signature complex is 2 and completes the proof of Lemma 3.1.4. This therefore completes the proof of the Hirzebruch signature theorem in general.

The de Rham complex with coeffiecients in $V$ is defined by the diagram:

$$
C^{\infty}\left(\Lambda^{\mathrm{e}, \mathrm{o}}(M) \otimes V\right) \xrightarrow{\nabla} C^{\infty}\left(T^{*} M \otimes \Lambda^{\mathrm{e}, \mathrm{o}}(M) \otimes V\right) \xrightarrow{c \otimes 1} C^{\infty}\left(\Lambda^{\mathrm{o}, \mathrm{e}}(M) \otimes V\right)
$$

and we shall denote the operator by $(d+\delta)_{V}^{e}$. The relations given by Lemma 3.4.2 give rise to relations among the local formulas:

$$
\begin{aligned}
& a_{n}\left(x,(d+\delta)_{V}^{\mathrm{e}}\right)=a_{n}\left(x, A_{(-1)^{m / 2}\left(\Delta^{+}-\Delta^{-}\right) \otimes V}^{+}\right) \\
& a_{n}\left(x,(d+\delta)_{V}^{+}\right)=a_{n}\left(x, A_{\left(\Delta^{+}-\Delta^{-}\right) \otimes V}^{+}\right)
\end{aligned}
$$

where $(d+\delta)_{V}^{+}$is the operator of the twisted signature complex. These relations are well defined regardless of whether or not $M$ admits a spin structrue on $T(M)$.

We deal first with the de Rham complex. Using Lemma 3.4.1 and the fact that the normalizing constant $c$ is 1 , we conclude:

$$
\left.a_{n}\left(x,(d+\delta)_{V}^{\mathrm{e}}\right)=\hat{A}^{\prime} \wedge \operatorname{ch}\left((-1)^{m / 2}\left(\Delta^{+}-\Delta^{-}\right)\right)\right) \wedge \operatorname{ch}(V) .
$$

Since $\operatorname{ch}\left((-1)^{m / 2}\left(\Delta^{+}{ }^{-} \Delta^{-}\right)\right)=e(M)$ is already a top dimensional form by Lemma 3.3.5(b), we conclude $a_{n}\left(x,(d+\delta)_{V}^{\mathrm{e}}\right)=(\operatorname{dim} V) e(M)$. This proves:
Theorem 3.4.3. Let $(d+\delta)_{V}^{\mathrm{e}}$ be the de Rham complex with coefficients in the bundle $V$ and let $a_{n}\left(x,(d+\delta)_{V}^{\mathrm{e}}\right)$ be the invariants of the heat equation. Then:
(a) $a_{n}\left(x,(d+\delta)_{V}^{e}\right)=0$ if $n<m$.
(b) $a_{m}\left(x,(d+\delta)_{V}^{\mathrm{e}}\right)=(\operatorname{dim} V) e(M)$ where $e(M)$ is the Euler form of $T(M)$. (c)

$$
\operatorname{index}\left((d+\delta)_{V}^{\mathrm{e}}\right)=(\operatorname{dim} V) \chi(M)=\int_{m}(\operatorname{dim} V) e(M)
$$

This shows that no information about $V$ (except its dimension) is obtained by considering the de Rham complex with coefficients in $V$ and it is for this reason we did not introduce this complex earlier. This gives a second proof of the Gauss-Bonnet theorem independent of the proof we gave earlier.

Next, we study the signature complex in order to compute $\hat{A}$. Using Lemma 3.4.1 with the bundle $V=1$, we conclude for $m=4 k$,

$$
a_{m}\left(x,(d+\delta)_{V}^{+}\right)=\left\{\hat{A}^{\prime} \wedge \operatorname{ch}(\Delta)\right\}_{m}
$$

Using Theorem 3.1.1 and Lemma 3.3.6, we compute therefore:

$$
L_{k}=\{\operatorname{ch}(\Delta) \wedge \hat{A}\}_{m}=\left\{\operatorname{ch}(\Delta) \wedge \hat{A}^{\prime}\right\}_{m}=a_{m}\left(x,(d+\delta)_{V}^{+}\right)
$$

so as the Chern character is formally invertible; $\hat{A}=\hat{A}^{\prime}$ is given by the generating function $z_{j} / \sinh \left(z_{j} / 2\right)$ by Lemma 3.3.7. We can now improve Lemma 3.4.1 and determine all the relevant normalizing constants.

Theorem 3.4.4. Let $T(M)$ admit a spin structure, then:
(a) $a_{n}^{\text {spin }}=0$ if $n<m$.
(b) $a_{m}^{n \text { spin }}(x, V)=\sum_{4 s+2 t=m} \hat{A}_{s} \wedge c h_{t}(V)$.
(c) $\operatorname{index}(V$, spin $)=\int_{M} a_{m}^{\text {spin }}(x, V)$.

### 3.5. The Riemann-Roch Theorem For Almost Complex Manifolds.

So far, we have discussed three of the four classical elliptic complexes, the de Rham, the signature, and the spin complexes. In this subsection, we define the Dolbeault complex for an almost complex manifold and relate it to the spin complex.

Let $M$ be a Riemannian manifold of dimension $m=2 n$. An almost complex structure on $M$ is a linear map $J: T(M) \rightarrow T(M)$ with $J^{2}=$ -1 . The Riemannian metric $G$ is unitary if $G(X, Y)=G(J X, J Y)$ for all $X, Y \in T(M)$; we can always construct unitary metrics by averaging over the action of $J$. Henceforth we assume $G$ is unitary. We extend $G$ to be Hermitian on $T(M) \otimes \mathbf{C}, T^{*}(M) \otimes \mathbf{C}$, and $\Lambda(M) \otimes \mathbf{C}$.

Since $J^{2}=-1$, we decompose $T(M) \otimes \mathbf{C}=T^{\prime}(M) \oplus T^{\prime \prime}(M)$ into the $\pm i$ eigenspaces of $J$. This direct sum is orthogonal with respect to the metric $G$. Let $\Lambda^{1,0}(M)$ and $\Lambda^{0,1}(M)$ be the dual spaces in $T^{*}(M) \otimes \mathbf{C}$ to $T^{\prime}$ and $T^{\prime \prime}$. We choose a local frame $\left\{e_{j}\right\}$ for $T(M)$ so that $J\left(e_{j}\right)=e_{j+n}$ for $1 \leq i \leq n$. Let $\left\{e^{j}\right\}$ be the corresponding dual frame for $T^{*}(M)$. Then:

$$
\begin{aligned}
T^{\prime}(M) & =\operatorname{span}_{\mathbf{C}}\left\{e_{j}-i e_{j+n}\right\}_{j=1}^{n} & T^{\prime \prime}(M) & =\operatorname{span}_{\mathbf{C}}\left\{e_{j}+i e_{j+n}\right\}_{j=1}^{n} \\
\Lambda^{1,0}(M) & =\operatorname{span}_{\mathbf{C}}\left\{e^{j}+i e^{j+n}\right\}_{j=1}^{n} & \Lambda^{0,1}(M) & =\operatorname{span}_{\mathbf{C}}\left\{e^{j}-i e^{j+n}\right\}_{j=1}^{n}
\end{aligned}
$$

These four vector bundles are all complex vector bundles over $M$. The metric gives rise to natural isomorphisms $T^{\prime}(M) \simeq \Lambda^{0,1}(M)$ and $T^{\prime \prime}(M)=$ $\Lambda^{1,0}(M)$. We will use this isomorphism to identify $e_{j}$ with $e^{j}$ for much of what follows.

If we forget the complex structure on $T^{\prime}(M)$, then the underlying real vector bundle is naturally isomorphic to $T(M)$. Complex multiplication by $i$ on $T^{\prime}(M)$ is equivalent to the endomorphism $J$ under this identification. Thus we may regard $J$ as giving a complex structure to $T(M)$.

The decomposition:

$$
T^{*}(M) \otimes \mathbf{C}=\Lambda^{1,0}(M) \oplus \Lambda^{0,1}(M)
$$

gives rise to a decomposition:

$$
\Lambda\left(T^{*} M\right) \otimes \mathbf{C}=\bigoplus_{p, q} \Lambda^{p, q}(M)
$$

for

$$
\Lambda^{p, q}(M)=\Lambda^{p}\left(\Lambda^{1,0}(M)\right) \otimes \Lambda^{q}\left(\Lambda^{0,1}(M)\right)
$$

Each of the bundles $\Lambda^{p, q}$ is a complex bundle over $M$ and this decomposition of $\Lambda\left(T^{*} M\right) \otimes \mathbf{C}$ is orthogonal. Henceforth we will denote these bundles by $T^{\prime}, T^{\prime \prime}$, and $\Lambda^{p, q}$ when no confusion will arise.

If $M$ is a holomorphic manifold, we let $z=\left(z_{1}, \ldots, z_{n}\right)$ be a local holomorphic coordinate chart. We expand $z_{j}=x_{j}+i y_{j}$ and define:

$$
\begin{aligned}
\frac{\partial}{\partial z_{j}} & =\frac{1}{2}\left(\frac{\partial}{\partial x_{j}}-i \frac{\partial}{\partial y_{j}}\right), & \frac{\partial}{\partial \bar{z}_{j}} & =\frac{1}{2}\left(\frac{\partial}{\partial x_{j}}+i \frac{\partial}{\partial y_{j}}\right) \\
d z_{j} & =d x_{j}+i d y_{j}, & d \bar{z}_{j} & =d x_{j}-i d y_{j} .
\end{aligned}
$$

We define:

$$
\partial(f)=\sum_{j} \frac{\partial f}{\partial z_{j}} d z_{j} \quad \text { and } \quad \bar{\partial}(f)=\sum_{j} \frac{\partial f}{\partial \bar{z}_{j}} d \bar{z}_{j} .
$$

The Cauchy-Riemann equations imply that a function $f$ is holomorphic if and only if $\bar{\partial}(f)=0$. If $w=\left(w_{1}, \ldots, w_{n}\right)$ is another holomorphic coordinate system, then:

$$
\begin{aligned}
\frac{\partial}{\partial w_{j}} & =\sum_{k} \frac{\partial z_{k}}{\partial w_{j}} \frac{\partial}{\partial z_{k}}, & \frac{\partial}{\partial \bar{w}_{j}} & =\sum_{k} \frac{\partial \bar{z}_{k}}{\partial \bar{w}_{j}} \frac{\partial}{\partial \bar{z}_{k}} \\
d w_{j} & =\sum_{k} \frac{\partial w_{j}}{\partial z_{k}} d z_{k}, & d \bar{w}_{j} & =\frac{\partial \bar{w}_{j}}{\partial \bar{z}_{k}} d \bar{z}_{k}
\end{aligned}
$$

We define:

$$
\begin{aligned}
T^{\prime}(M) & =\operatorname{span}\left\{\frac{\partial}{\partial z_{j}}\right\}_{j=1}^{n}, & T^{\prime \prime}(M) & =\operatorname{span}\left\{\frac{\partial}{\partial \bar{z}_{j}}\right\}_{j=1}^{n} \\
\Lambda^{1,0}(M) & =\operatorname{span}\left\{d z_{j}\right\}_{j=1}^{n}, & \Lambda^{0,1}(M) & =\operatorname{span}\left\{d \bar{z}_{j}\right\}_{j=1}^{n}
\end{aligned}
$$

then these complex bundles are invariantly defined independent of the choice of the coordinate system. We also note

$$
\partial: C^{\infty}(M) \rightarrow C^{\infty}\left(\Lambda^{1,0}(M)\right) \quad \text { and } \quad \bar{\partial}: C^{\infty}(M) \rightarrow C^{\infty}\left(\Lambda^{0,1}(M)\right)
$$

are invariantly defined and decompose $d=\partial+\bar{\partial}$.
There is a natural isomorphism of $T^{\prime}(M)$ with $T(M)$ as real bundles in this example and we let $J$ be complex multiplication by $i$ under this isomorphism. Equivalently:

$$
J\left(\frac{\partial}{\partial x_{j}}\right)=\left(\frac{\partial}{\partial y_{j}}\right) \quad \text { and } \quad J\left(\frac{\partial}{\partial y_{j}}\right)=-\frac{\partial}{\partial x_{j}} .
$$

$T^{\prime}(M)=T_{c}(M)$ is the complex tangent bundle in this example; we shall reserve the notation $T_{c}(M)$ for the holomorphic case.

Not every almost complex structure arises from a complex structure; there is an integrability condition. If $J$ is an almost complex structure, we
decompose the action of exterior differentiation $d$ on $C^{\infty}(\Lambda)$ with respect to the bigrading $(p, q)$ to define:

$$
\partial: C^{\infty}\left(\Lambda^{p, q}\right) \rightarrow C^{\infty}\left(\Lambda^{p+1, q}\right) \quad \text { and } \quad \bar{\partial}: C^{\infty}\left(\Lambda^{p, q}\right) \rightarrow C^{\infty}\left(\Lambda^{p, q+1}\right)
$$

Theorem 3.5.1 (Nirenberg-Neulander). The following are equivalent and define the notion of an integrable almost complex structure:
(a) The almost complex structure $J$ arises from a holomorphic structure on $M$.
(b) $d=\partial+\bar{\partial}$.
(c) $\bar{\partial} \bar{\partial}=0$.
(d) $T^{\prime}(M)$ is integrable-i.e., given $X, Y \in C^{\infty}\left(T^{\prime}(M)\right)$, then the Lie bracket $[X, Y] \in C^{\infty}\left(T^{\prime}(M)\right)$ where we extend [ , ] to complex vector fields in the obvious fashion.
Proof: This is a fairly deep result and we shall not give complete details. It is worth, however, giving a partial proof of some of the implications to illustrate the concepts we will be working with. Suppose first $M$ is holomorphic and let $\left\{z_{j}\right\}$ be local holomorphic coordinates on $M$. Define:

$$
d z^{I}=d z_{i_{1}} \wedge \cdots \wedge d z_{i_{p}} \quad \text { and } \quad d \bar{z}^{J}=d \bar{z}_{j_{1}} \wedge \cdots \wedge d \bar{z}_{j_{q}}
$$

then the collection $\left\{d z^{I} \wedge d \bar{z}^{J}\right\}$ gives a local frame for $\Lambda^{p, q}$ which is closed. If $\omega \in C^{\infty}\left(\Lambda^{p, q}\right)$, we decompose $\omega=\sum f_{I, J} d z^{I} \wedge d \bar{z}^{J}$ and compute:

$$
d \omega=d\left\{\sum f_{I, J} d z^{I} \wedge d \bar{z}^{J}\right\}=\sum d f_{I, J} \wedge d z^{I} \wedge d \bar{z}^{J}
$$

On functions, we decompose $d=\partial+\bar{\partial}$. Thus $d \omega \in C^{\infty}\left(\Lambda^{p+1, q} \oplus \Lambda^{p, q+1}\right)$ has no other components. Therefore $d \omega=\partial \omega+\bar{\partial} \omega$ so (a) implies (b).

We use the identity $d^{2}=0$ to compute $(\partial+\bar{\partial})^{2}=\left(\partial^{2}\right)+(\partial \bar{\partial}+\bar{\partial} \partial)+$ $\left(\bar{\partial}^{2}\right)=0$. Using the bigrading and decomposing this we conclude $\left(\partial^{2}\right)=$ $(\partial \bar{\partial}+\bar{\partial} \partial)=\left(\bar{\partial}^{2}\right)=0$ so (b) implies (c). Conversely, suppose that $\bar{\partial} \bar{\partial}=0$ on $C^{\infty}(M)$. We must show $d: C^{\infty}\left(\Lambda^{p, q}\right) \rightarrow C^{\infty}\left(\Lambda^{p+1, q} \oplus \Lambda^{p, q+1}\right)$ has no other components. Let $\left\{e_{j}\right\}$ be a local frame for $\Lambda^{1,0}$ and let $\left\{\bar{e}_{j}\right\}$ be the corresponding local frame for $\Lambda^{0,1}$. We decompose:

$$
\begin{aligned}
d e_{j}=\partial e_{j}+\bar{\partial} e_{j}+A_{j} & \text { for } A_{j} \in \Lambda^{0,2} \\
d \bar{e}_{j}=\bar{\partial} \bar{e}_{j}+\partial \bar{e}_{j}+\bar{A}_{j} & \text { for } \bar{A}_{j} \in \Lambda^{2,0}
\end{aligned}
$$

Then we compute:

$$
\begin{aligned}
d\left(\sum f_{j} e_{j}\right) & =\sum d f_{j} \wedge e_{j}+f_{j} d e_{j} \\
& =\sum\left\{\partial f_{j} \wedge e_{j}+\bar{\partial} f_{j} \wedge e_{j}+f_{j} \partial e_{j}+f_{j} \bar{\partial} e_{j}+f_{j} A_{j}\right\} \\
& =\partial\left(\sum f_{j} e_{j}\right)+\bar{\partial}\left(\sum f_{j} e_{j}\right)+\sum f_{j} A_{j}
\end{aligned}
$$

Similarly

$$
d\left(\sum f_{j} \bar{e}_{j}\right)=\partial\left(\sum f_{j} \bar{e}_{j}\right)+\bar{\partial}\left(\sum f_{j} \bar{e}_{j}\right)+\sum f_{j} \bar{A}_{j}
$$

If $A$ and $\bar{A}$ denote the $0^{\text {th }}$ order operators mapping $\Lambda^{1,0} \rightarrow \Lambda^{0,2}$ and $\Lambda^{0,1} \rightarrow$ $\Lambda^{2,0}$ then we compute:

$$
\begin{array}{ll}
d(\omega)=\partial(\omega)+\bar{\partial}(\omega)+A \omega & \text { for } \omega \in C^{\infty}\left(\Lambda^{1,0}\right) \\
d(\bar{\omega})=\partial(\bar{\omega})+\bar{\partial}(\bar{\omega})+\bar{A} \bar{\omega} & \text { for } \bar{\omega} \in C^{\infty}\left(\Lambda^{0,1}\right) .
\end{array}
$$

Let $f \in C^{\infty}(M)$ and compute:

$$
0=d^{2} f=d(\partial f+\bar{\partial} f)=\left(\partial^{2}+\bar{A} \bar{\partial}\right) f+(\partial \bar{\partial}+\bar{\partial} \partial) f+\left(\bar{\partial}^{2}+A \partial\right) f
$$

where we have decomposed the sum using the bigrading. This implies that $\left(\bar{\partial}^{2}+A \partial\right) f=0$ so $A \partial f=0$. Since $\{\partial f\}$ spans $\Lambda^{1,0}$, this implies $A=\bar{A}=0$ since $A$ is a $0^{\text {th }}$ order operator. Thus $d e_{j}=\partial e_{j}+\bar{\partial} e_{j}$ and $d \bar{e}_{j}=\partial \bar{e}_{j}+\bar{\partial} \bar{e}_{j}$. We compute $d\left(e^{I} \wedge \bar{e}^{J}\right)$ has only $(p+1, q)$ and $(p, q+1)$ components so $d=\partial+\bar{\partial}$. Thus (b) and (c) are equivalent.

It is immediate from the definition that $X \in T^{\prime}(M)$ if and only if $\omega(X)=$ 0 for all $\omega \in \Lambda^{0,1}$. If $X, Y \in C^{\infty}\left(T^{*} M\right)$, then Cartan's identity implies:

$$
d \omega(X, Y)=X(\omega Y)-Y(\omega X)-\omega([X, Y])=-\omega([X, Y])
$$

If (b) is true then $d \omega$ has no component in $\Lambda^{2,0}$ so $d \omega(X, Y)=0$ which implies $\omega([X, Y])=0$ which implies $[X, Y] \in C^{\infty}\left(T^{\prime} M\right)$ which implies (d). Conversly, if (d) is true, then $d \omega(X, Y)=0$ so $d \omega$ has no component in $\Lambda^{2,0}$ so $d=\partial+\bar{\partial}$ on $\Lambda^{0,1}$. By taking conjugates, $d=\partial+\bar{\partial}$ on $\Lambda^{1,0}$ as well which implies as noted above that $d=\partial+\bar{\partial}$ in general which implies (c).

We have proved that (b)-(d) are equivalent and that (a) implies (b). The hard part of the theorem is showing (b) implies (a). We shall not give this proof as it is quite lengthy and as we shall not need this implication of the theorem.

As part of the previous proof, we computed that $d-(\partial+\bar{\partial})$ is a $0^{\text {th }}$ order operator (which vanishes if and only if $M$ is holomorphic). We now compute the symbol of both $\partial$ and $\bar{\partial}$. We use the metric to identify $T(M)=T^{*}(M)$. Let $\left\{e_{j}\right\}$ be a local orthonormal frame for $T(M)$ such that $J\left(e_{j}\right)=e_{j+n}$ for $1 \leq j \leq n$. We extend ext and int to be complex linear maps from $T^{*}(M) \otimes \mathbf{C} \rightarrow \operatorname{END}\left(\Lambda\left(T^{*} M\right) \otimes \mathbf{C}\right)$.

Lemma 3.5.2. Let $\partial$ and $\bar{\partial}$ be defined as before and let $\delta^{\prime}$ and $\delta^{\prime \prime}$ be the formal adjoints. Then:
(a)

$$
\begin{aligned}
\partial: C^{\infty}\left(\Lambda^{p, q}\right) & \rightarrow C^{\infty}\left(\Lambda^{p+1, q}\right) \\
\text { and } & \sigma_{L}(\partial)(x, \xi)=\frac{i}{2} \sum_{j \leq n}\left(\xi_{j}-i \xi_{j+n}\right) \operatorname{ext}\left(e_{j}+i e_{j+n}\right) \\
\bar{\partial}: C^{\infty}\left(\Lambda^{p, q}\right) & \rightarrow C^{\infty}\left(\Lambda^{p, q+1}\right) \\
\text { and } & \sigma_{L}(\bar{\partial})(x, \xi)=\frac{i}{2} \sum_{j \leq n}\left(\xi_{j}+i \xi_{j+n}\right) \operatorname{ext}\left(e_{j}-i e_{j+n}\right) \\
\delta^{\prime}: C^{\infty}\left(\Lambda^{p, q}\right) & \rightarrow C^{\infty}\left(\Lambda^{p-1, q}\right) \\
\text { and } & \sigma_{L}\left(\delta^{\prime}\right)(x, \xi)=-\frac{i}{2} \sum_{j \leq n}\left(\xi_{j}+i \xi_{j+n}\right) \operatorname{int}\left(e_{j}-i e_{j+n}\right) \\
\delta^{\prime \prime}: C^{\infty}\left(\Lambda^{p, q}\right) & \rightarrow C^{\infty}\left(\Lambda^{p, q-1}\right) \\
\text { and } & \sigma_{L}\left(\delta^{\prime \prime}\right)(x, \xi)=-\frac{i}{2} \sum_{j \leq n}\left(\xi_{j}-i \xi_{j+n}\right) \operatorname{int}\left(e_{j}+i e_{j+n}\right) .
\end{aligned}
$$

(b) If $\Delta_{c}^{\prime}=\left(\partial+\delta^{\prime}\right)^{2}$ and $\Delta_{c}^{\prime \prime}=\left(\bar{\partial}+\delta^{\prime \prime}\right)^{2}$ then these are elliptic on $\Lambda$ with $\sigma_{L}\left(\Delta_{c}^{\prime}\right)=\sigma_{l}\left(\Delta_{c}^{\prime \prime}\right)=\frac{1}{2}|\xi|^{2} I$.
Proof: We know $\sigma_{L}(d)(x, \xi)=i \sum_{j} \xi_{j} \operatorname{ext}\left(e_{j}\right)$. We define

$$
A(\xi)=\frac{1}{2} \sum_{j \leq n}\left(\xi_{j+n}\right) \operatorname{ext}\left(e_{j}+i e_{j+n}\right)
$$

then $A(\xi): \Lambda^{p, q} \rightarrow \Lambda^{p+1, q}$ and $\bar{A}(\xi): \Lambda^{p, q} \rightarrow \Lambda^{p, q+1}$. Since $i A(\xi)+i \bar{A}(\xi)=$ $\sigma_{L}(d)$, we conclude that $i A$ and $i \bar{A}$ represent the decomposition of $\sigma_{L}(d)$ under the bigrading and thus define the symbols of $\partial$ and $\bar{\partial}$. The symbol of the adjoint is the adjoint of the symbol and this proves (a). (b) is an immediate consequence of (a).

If $M$ is holomorphic, the Dolbeault complex is the complex $\left\{\bar{\partial}, \Lambda^{0, q}\right\}$ and the index of this complex is called the arithmetic genus of $M$. If $M$ is not holomorphic, but only has an almost complex structure, then $\bar{\partial}^{2} \neq 0$ so we can not define the arithmetic genus in this way. Instead, we use a trick called "rolling up" the complex. We define:

$$
\Lambda^{0,+}=\bigoplus_{q} \Lambda^{0,2 q} \quad \text { and } \quad \Lambda^{0,-}=\bigoplus_{q} \Lambda^{0,2 q+1}
$$

to define a $\mathbf{Z}_{2}$ grading on the Dolbeault bundles. (These are also often denoted by $\Lambda^{0, \text { even }}$ and $\Lambda^{0, \text { odd }}$ in the literature). We consider the two term elliptic complex:

$$
\left(\bar{\partial}+\delta^{\prime \prime}\right)_{ \pm}: C^{\infty}\left(\Lambda^{0, \pm}\right) \rightarrow C^{\infty}\left(\Lambda^{0, \mp}\right)
$$

and define the arithmetic genus of $M$ to be the index of $\left(\bar{\partial}+\delta^{\prime \prime}\right)_{+}$. The adjoint of $\left(\bar{\partial}+\delta^{\prime \prime}\right)_{+}$is $\left(\bar{\partial}+\delta^{\prime \prime}\right)_{-}$and the associated Laplacian is just $\Delta_{c}^{\prime \prime}$ restricted to $\Lambda^{0, \pm}$ so this is an elliptic complex. If $M$ is holomorphic, then the index of this elliptic complex is equal to the index of the complex ( $\bar{\partial}, \Lambda^{0, q}$ ) by the Hodge decomposition theorem.

We define:

$$
\begin{aligned}
& c^{\prime}(\xi)=(\sqrt{2})^{-1} \sum_{j \leq n}\left\{\left(\xi_{j}+i \xi_{j+n}\right) \operatorname{ext}\left(e_{j}-i e_{j+n}\right)\right. \\
&\left.-\left(\xi_{j}-i \xi_{j+n}\right) \operatorname{int}\left(e_{j}+i e_{j+n}\right)\right\}
\end{aligned}
$$

then it is immediate that:

$$
c^{\prime}(\xi) c^{\prime}(\xi)=-|\xi|^{2} \quad \text { and } \quad \sigma_{L}\left(\bar{\partial}+\delta^{\prime \prime}\right)=i c^{\prime}(\xi) / \sqrt{2}
$$

Let $\nabla$ be a connection on $\Lambda^{0, \pm}$, then we define the operator $A^{ \pm}(\nabla)$ by the diagram:

$$
A^{ \pm}(\nabla): C^{\infty}\left(\Lambda^{0, \pm}\right) \xrightarrow{\nabla} C^{\infty}\left(T^{*} M \otimes \Lambda^{0, \pm}\right) \xrightarrow{c^{\prime} / \sqrt{2}} C^{\infty}\left(\Lambda^{0, \mp}\right)
$$

This will have the same leading symbol as $\left(\bar{\partial}+\delta^{\prime \prime}\right)_{ \pm}$. There exists a unique connection so $A^{ \pm}(\nabla)=\left(\bar{\partial}+\delta^{\prime \prime}\right)_{ \pm}$but as the index is constant under lower order perturbations, we shall not need this fact as the index of $A^{+}(\nabla)=\operatorname{index}\left(\bar{\partial}+\delta^{\prime \prime}\right)_{+}$for any $\nabla$. We shall return to this point in the next subsection.

We now let $V$ be an arbitrary coefficient bundle with a connection $\nabla$ and define the Dolbeault complex with coefficients in $V$ using the diagram:

$$
A^{ \pm}: C^{\infty}\left(\Lambda^{0, \pm} \otimes V\right) \xrightarrow{\nabla} C^{\infty}\left(T^{*} M \otimes \Lambda^{0, \pm} \otimes V\right) \xrightarrow{c^{\prime} / \sqrt{2} \otimes 1} C^{\infty}\left(\Lambda^{0, \mp} \otimes V\right)
$$

and we define index ( $V$, Dolbeault) to be the index of this elliptic complex. The index is independent of the connections chosen, of the fiber metrics chosen, and is constant under perturbations of the almost complex structure.

If $M$ is holomorphic and if $V$ is a holomorphic vector bundle, then we can extend $\bar{\partial}: C^{\infty}\left(\Lambda^{0, q} \otimes V\right) \rightarrow C^{\infty}\left(\Lambda^{0, q+1} \otimes V\right)$ with $\bar{\partial} \bar{\partial}=0$. Exactly as was true for the arithmetic genus, the index of this elliptic complex is equal to the index of the rolled up elliptic complex so our definitions generalize the usual definitions from the holomorphic category to the almost complex category.

We will compute a formula for index( $V$, Dolbeault) using the spin complex. There is a natural inclusion from $\mathrm{U}\left(\frac{m}{2}\right)=\mathrm{U}(n)$ into $\mathrm{SO}(m)$, but this does not lift in general to $\operatorname{SPIN}(m)$. We saw earlier that $T\left(\mathbf{C} P_{k}\right)$
does not admit a spin structure if $k$ is even, even though it does admit a unitary structure. We define $\operatorname{SPIN}_{c}(m)=\operatorname{SPIN}(m) \times S^{1} / \mathbf{Z}_{2}$ where we choose the $\mathbf{Z}_{2}$ identification $(g, \lambda)=(-g,-\lambda)$. There is a natural map $\rho_{c}: \operatorname{SPIN}_{c}(m) \rightarrow \mathrm{SO}(m) \times S^{1}$ induced by the map which sends $(g, \lambda) \mapsto$ $\left(\rho(g), \lambda^{2}\right)$. This map is a $\mathbf{Z}_{2}$ double cover and is a group homomorphism.

There is a natural map of $\mathrm{U}\left(\frac{m}{2}\right)$ into $\mathrm{SO}(m) \times S^{1}$ defined by sending $g \mapsto(g, \operatorname{det}(g))$. The interesting thing is that this inclusion does lift; we can define $f: \mathrm{U}\left(\frac{m}{2}\right) \rightarrow \mathrm{SPIN}_{c}(m)$ so the following diagram commutes:


We define the lifting as follows. For $U \in \mathrm{U}\left(\frac{m}{2}\right)$, we choose a unitary basis $\left\{e_{j}\right\}_{j=1}^{n}$ so that $U\left(e_{j}\right)=\lambda_{j} e_{j}$. We define $e_{j+n}=i e_{j}$ so $\left\{e_{j}\right\}_{j=1}^{2 n}$ is an orthogonal basis for $\mathbf{R}^{m}=\mathbf{C}^{n}$. Express $\lambda_{j}=e^{i \theta_{j}}$ and define:

$$
f(U)=\prod_{j=1}^{n}\left\{\cos \left(\theta_{j} / 2\right)+\sin \left(\theta_{j} / 2\right) e_{j} e_{j+n}\right\} \times \prod_{j=1}^{n} e^{i \theta_{j} / 2} \in \operatorname{SPIN}_{c}(m)
$$

We note first that $e_{j} e_{j+n}$ is an invariant of the one-dimensional complex subspace spanned by $e_{j}$ and does not change if we replace $e_{j}$ by $z e_{j}$ for $|z|=1$. Since all the factors commute, the order in which the eigenvalues is taken does not affect the product. If there is a multiple eigenvalue, this product is independent of the particular basis which is chosen. Finally, if we replace $\theta_{j}$ by $\theta_{j}+2 \pi$ then both the first product and the second product change sign. Since $(g, \lambda)=(-g,-\lambda)$ in $\operatorname{SPIN}_{c}(m)$, this element is invariantly defined. It is clear that $f(I)=I$ and that $f$ is continuous. It is easily verified that $\rho_{c} f(U)=i(U) \times \operatorname{det}(U)$ where $i(U)$ denotes the matrix $U$ viewed as an element of $\operatorname{SO}(m)$ where we have forgotten the complex structure on $\mathbf{C}^{n}$. This proves that $f$ is a group homomorphism near the identity and consequently $f$ is a group homomorphism in general. $\rho_{c}$ is a covering projection.

The cannonical bundle $K$ is given by:

$$
K=\Lambda^{n, 0}=\Lambda^{n}\left(T^{\prime} M\right)^{*}
$$

so $K^{*}=\Lambda^{n}\left(T^{\prime} M\right)$ is the bundle with transition functions $\operatorname{det}\left(U_{\alpha \beta}\right)$ where the $U_{\alpha \beta}$ are the transitions for $T^{\prime} M$. If we have a spin structure on $M$, we can split:

$$
P_{\mathrm{SPIN}_{c}}=P_{\mathrm{SPIN}} \times P_{S^{1}} / \mathbf{Z}_{2}
$$

where $P_{S^{1}}$ represents a line bundle $L_{1}$ over $M$. From this description, it is clear $K^{*}=L_{1} \otimes L_{1}$. Conversely, if we can take the square root of $K$ (or equivalently of $K^{*}$ ), then $M$ admits a spin structure so the obstruction to constructing a spin structure on an almost complex manifold is the obstruction to finding a square root of the cannonical bundle.

Let $V=\mathbf{C}^{n}=\mathrm{R}^{m}$ with the natural structures. We extend $\Delta^{ \pm}$to representations $\Delta^{ \pm}$of $\operatorname{SPIN}_{c}(m)$ in the natural way. We relate this representation to the Dolbeault representation as follows:

Lemma 3.5.3. There is a natural isomorphism between $\Lambda^{0, \pm}$ and $\Delta_{c}^{ \pm}$ which defines an equivalence of these two representations of $\mathrm{U}\left(\frac{m}{2}\right)$. Under this isomorphism, the action of $V$ by Clifford multiplication on the left is preserved.

Proof: Let $\left\{e_{j}\right\}$ be an orthonormal basis for $\mathbf{R}^{m}$ with $J e_{j}=e_{j+n}$ for $1 \leq j \leq n$. Define:
$\alpha_{j}=i e_{j} e_{j+n}, \quad \beta_{j}=e_{j}+i e_{j+n}, \quad \bar{\beta}_{j}=e_{j}-i e_{j+n} \quad$ for $1 \leq j \leq n$.
We compute:

$$
\beta_{j} \alpha_{j}=-\beta_{j} \quad \text { and } \quad \beta_{j} \alpha_{k}=\alpha_{k} \beta_{j} \quad \text { for } k \neq j
$$

We define $\gamma=\beta_{1} \ldots \beta_{n}$ then $\gamma$ spans $\Lambda^{n, 0}$ and $\gamma \alpha_{j}=-\gamma, 1 \leq j \leq n$. We define:

$$
\Lambda^{0, *}=\bigoplus_{q} \Lambda^{0, q} \quad \text { and } \quad \Lambda^{n, *}=\bigoplus_{q} \Lambda^{n, q}=\Lambda^{0, *} \gamma
$$

Since $\operatorname{dim}\left(\Lambda^{0, *}\right)=2^{n}$, we conclude that $\Lambda^{0, *} \gamma$ is the simultaneous -1 eigenspace for all the $\alpha_{j}$ and that therefore:

$$
\Delta_{c}=\Lambda^{0, *} \gamma
$$

as a left representation space for $\operatorname{SPIN}(m)$. Again, we compute:

$$
\alpha_{j} \beta_{j}=\beta_{j}, \quad \alpha_{j} \bar{\beta}_{j}=-\bar{\beta}_{j}, \quad \alpha_{j} \beta_{k}=\beta_{k} \alpha_{j}, \quad \alpha_{j} \bar{\beta}_{k}=\bar{\beta}_{k} \alpha_{j} \quad \text { for } j \neq k
$$

so that if $x \in \Lambda^{0, q} \gamma$ then since $\tau=\alpha_{1} \ldots \alpha_{n}$ we have $\tau x=(-1)^{q} x$ so that

$$
\Delta_{c}^{ \pm}=\Lambda^{0, \pm} \gamma
$$

We now study the induced representation of $\mathrm{U}(n)$. Let

$$
g=\left(\cos (\theta / 2)+\sin (\theta / 2) e_{1} J e_{1}\right) \cdot e^{i \theta / 2} \in \operatorname{SPIN}_{c}(m)
$$

We compute:

$$
\begin{aligned}
& g \beta_{1}=\{\cos (\theta / 2)-i \sin (\theta / 2)\} e^{i \theta / 2} \beta_{1}=\beta_{1} \\
& g \bar{\beta}_{1}=\{\cos (\theta / 2)+i \sin (\theta / 2)\} e^{i \theta / 2} \bar{\beta}_{1}=e^{i \theta} \bar{\beta}_{1} \\
& g \beta_{k}=\beta_{k} g \quad \text { and } \quad g \bar{\beta}_{k}=\bar{\beta}_{k} g \quad \text { for } k>1
\end{aligned}
$$

Consequently:

$$
g \bar{\beta}_{j} \gamma= \begin{cases}\bar{\beta}_{j} \gamma & \text { if } j_{1}>1 \\ e^{i \theta} \bar{\beta}_{j} \gamma & \text { if } j_{1}=1\end{cases}
$$

A similar computation goes for all the other indices and thus we compute that if $U e_{j}=e^{i \theta_{j}} e_{j}$ is unitary that:

$$
f(U) \bar{\beta}_{J} \gamma=e^{i \theta_{j_{1}}} \ldots e^{i \theta_{j_{q}}} \bar{\beta}_{J} \gamma
$$

which is of course the natural represenation of $\mathrm{U}(n)$ on $\Lambda^{0, q}$.
Finally we compare the two actions of $V$ by Clifford multiplication. We assume without loss of generality that $\xi=(1,0, \ldots, 0)$ so we must study Clifford multiplication by $e_{1}$ on $\Delta_{c}$ and

$$
\left\{\operatorname{ext}\left(e_{1}-i e_{1+n}\right)-\operatorname{int}\left(e_{1}+i e_{1+n}\right)\right\} / \sqrt{2}=c^{\prime}\left(e_{1}\right)
$$

on $\Lambda^{0, *}$. We compute:

$$
\begin{aligned}
e_{1} \beta_{1} & =-1+i e_{1} e_{2}=\left(e_{1}-i e_{2}\right)\left(e_{1}+i e_{2}\right) / 2=\bar{\beta}_{1} \beta_{1} / 2, \\
e_{1} \bar{\beta}_{1} \beta_{1} & =-2 \beta_{1}, \\
e_{1} \bar{\beta}_{k} & =-\bar{\beta}_{k} e_{1} \quad \text { for } k>1 .
\end{aligned}
$$

From this it follows immediately that:

$$
e_{1} \bar{\beta}_{J} \gamma=(-1)^{q} \begin{cases}\bar{\beta}_{1} \bar{\beta}_{J} \gamma / 2 & \text { if } j_{1}>1 \\ 2 \bar{\beta}_{J^{\prime}} \gamma & \text { where } J^{\prime}=\left\{j_{2}, \ldots, j_{q}\right\} \text { if } j_{1}=1\end{cases}
$$

Similarly, we compute:

$$
c^{\prime}\left(e_{1}\right) \bar{\beta}_{J}= \begin{cases}\bar{\beta}_{1} \bar{\beta}_{J} / \sqrt{2} & \text { if } j_{1}>1 \\ -\sqrt{2} \bar{\beta}_{J^{\prime}} & \text { for } J^{\prime}=\left\{j_{2}, \ldots, j_{q}\right\} \text { if } j_{1}=1\end{cases}
$$

From these equations, it is immediate that if we define $T\left(\bar{\beta}_{J}\right)=\bar{\beta}_{J} \gamma$, then $T$ will not preserve Clifford multiplication. We let $a(q)$ be a sequence of non-zero constants and define $T: \Lambda^{0, \pm} \rightarrow \Delta_{c}^{ \pm}$by:

$$
T\left(\bar{\beta}_{J}\right)=a(|J|) \bar{\beta}_{J} \gamma
$$

Since the spaces $\Lambda^{0, q}$ are $\mathrm{U}(n)$ invariant, this defines an equivalence between these two representations of $\mathrm{U}(n)$. $T$ will induce an equivalence between Clifford multiplication if and only if we have the relations:

$$
(-1)^{q} a(q) / 2=a(q+1) / \sqrt{2} \quad \text { and } \quad \sqrt{2} a(q-1)=(-1)^{q} \cdot 2 \cdot a(q)
$$

These give rise to the inductive relations:

$$
a(q+1)=(-1)^{q} a(q) / \sqrt{2} \quad \text { and } \quad a(q)=(-1)^{q-1} a(q-1) / \sqrt{2} .
$$

These relations are consistent and we set $a(q)=(\sqrt{2})^{-q}(-1)^{q(q-1) / 2}$ to define the equivalence $T$ and complete the proof of the Lemma.

From this lemma, we conclude:
Lemma 3.5.4. There is a natural isomorphism of elliptic complexes $\left(\Delta_{c}^{+}-\Delta_{c}^{-}\right) \otimes V \simeq\left(\Lambda^{0,+}-\Lambda^{0,-}\right) \otimes V$ which takes the operator of the $\mathrm{SPIN}_{c}$ complex to an operator which has the same leading symbol as the Dolbeault complex (and thus has the same index). Furthermore, we can represent the $\mathrm{SPIN}_{c}$ complex locally in terms of the SPIN complex in the form: $\left(\Delta_{c}^{+}-\Delta_{c}^{-}\right) \otimes V=\left(\Delta^{+}-\Delta^{-}\right) \otimes L_{1} \otimes V$ where $L_{1}$ is a local square root of $\Lambda^{n}\left(T^{\prime} M\right)$.

We use this sequence of isomorphisms to define an operator on the Dolbeault complex with the same leading symbol as the operator $\left(\bar{\partial}+\partial^{\prime \prime}\right)$ which is locally isomorphic to the natural operator of the SPIN complex. The $\mathbf{Z}_{2}$ ambiguity in the definition of $L_{1}$ does not affect this construction. (This is equivalent to choosing an appropriate connection called the spin connection on $\Lambda^{0, \pm}$.) We can compute index ( $V$, Dolbeault) using this operator. The local invariants of the heat equation for this operator are the local invariants of the twisted spin complex and therefore arguing exactly as we did for the signature complex, we compute:

$$
\operatorname{index}(V, \text { Dolbeault })=\int_{M} \hat{A}(T M) \wedge \operatorname{ch}\left(L_{1}\right) \wedge \operatorname{ch}(V)
$$

where $\operatorname{ch}\left(L_{1}\right)$ is to be understood as a complex characteristic class of $T^{\prime}(M)$.

Theorem 3.5.5 (Riemann-Roch). Let $\operatorname{Td}\left(T^{\prime} M\right)$ be the Todd class defined earlier by the generating function $\operatorname{Td}(A)=\prod_{\nu} x_{\nu} /\left(1-e^{-x_{\nu}}\right)$. Then

$$
\operatorname{index}(V, \text { Dolbeault })=\int_{M} T d\left(T^{\prime} M\right) \wedge c h(V)
$$

Proof: We must simply identify $\hat{A}(T M) \wedge c h\left(L_{1}\right)$ with $T d\left(T^{\prime} M\right)$. We perform a computation in characteristic classes using the splitting principal.

We formally decompose $T^{\prime} M=\tilde{L}_{1} \oplus \cdots \oplus \tilde{L}_{n}$ as the direct sum of line bundles. Then $\Lambda^{n}\left(T^{\prime} M\right)=\tilde{L}_{1} \otimes \cdots \otimes \tilde{L}_{n}$ so $c_{1}\left(\Lambda^{n} T^{\prime} M\right)=c_{1}\left(\tilde{L}_{1}\right)+$ $\cdots+c_{1}\left(\tilde{L}_{n}\right)=x_{1}+\cdots+x_{n}$. Since $L_{1} \otimes L_{1}=\Lambda^{n} T^{\prime} M$, we conclude $c_{1}\left(L_{1}\right)=\frac{1}{2}\left(x_{1}+\cdots+x_{n}\right)$ so that:

$$
\operatorname{ch}\left(L_{1}\right)=\prod_{\nu} e^{x_{\nu} / 2}
$$

Therefore:

$$
\begin{aligned}
\hat{A}(M) \wedge \operatorname{ch}\left(L_{1}\right) & =\prod_{\nu} x_{\nu} \frac{e^{x_{\nu} / 2}}{e^{x_{\nu} / 2}-e^{-x_{\nu} / 2}} \\
& =\prod_{\nu} \frac{x_{\nu}}{1-e^{-x_{\nu}}}=\operatorname{Td}\left(T^{\prime} M\right)
\end{aligned}
$$

Of course, this procedure is only valid if the given bundle does in fact split as the direct sum of line bundles. We can make this procedure a correct way of calculating characteristic classes by using flag manifolds or by actually calculating on the group representaton and using the fact that the diagonizable matrices are dense.

It is worth giving another proof of the Riemann-Roch formula to ensure that we have not made a mistake of sign somewhere in all our calculations. Using exactly the same multiplicative considerations as we used earlier and using the qualitative form of the formala for index ( $V$, Dolbeault) given by the $\mathrm{SPIN}_{c}$ complex, it is immediate that there is some formula of the form:

$$
\operatorname{index}(V, \text { Dolbeault })=\int_{N} \sum_{s+t=n} T d_{s}^{\prime}\left(T^{\prime} M\right) \wedge c^{t} c h_{t}(V)
$$

where $c$ is some universal constant to be determined and where $T d_{s}^{\prime}$ is some characteristic form of $T^{\prime} M$.

If $M=\mathbf{C} P_{j}$, then we shall show in Lemma 3.6.8 that the arithmetic genus of $\mathbf{C} P_{j}$ is 1 . Using the multiplicative property of the Dolbeault complex, we conclude the arithmetic genus of $\mathbf{C} P_{j_{1}} \times \cdots \times \mathbf{C} P_{j_{k}}$ is 1 as well. If $M_{\rho}^{c}$ are the manifolds of Lemma 2.3.4 then if we take $V=1$ we conclude:

$$
\begin{aligned}
1 & =\operatorname{index}(1, \text { Dolbeault })=\text { arithmetic genus of } \mathbf{C} P_{j_{1}} \times \cdots \times \mathbf{C} P_{j_{k}} \\
& =\int_{M_{\rho}^{c}} T d^{\prime}\left(T^{\prime} M\right)
\end{aligned}
$$

We verified in Lemma 2.3.5 that $T d\left(T^{\prime} M\right)$ also has this property so by the uniqueness assertion of Lemma 2.3.4 we conclude $T d=T d^{\prime}$. We take $m=2$ and decompose:

$$
\Lambda^{0} \oplus \Lambda^{2}=\Lambda^{0,0} \oplus \Lambda^{1,1} \quad \text { and } \quad \Lambda^{1}=\Lambda^{1,0} \oplus \Lambda^{0,1}
$$

so that

$$
\Lambda^{\mathrm{e}}-\Lambda^{\circ}=\left(\Lambda^{0,0}-\Lambda^{0,1}\right) \otimes\left(\Lambda^{0,0}-\Lambda^{1,0}\right)
$$

In dimension 2 (and more generally if $M$ is Kaehler), it is an easy exercise to compute that $\Delta=2\left(\bar{\partial} \delta^{\prime \prime}+\delta^{\prime \prime} \bar{\partial}\right)$ so that the harmonic spaces are the same therefore:

$$
\begin{aligned}
\chi(M) & =\operatorname{index}\left(\Lambda^{0,0}, \text { Dolbeault }\right)-\operatorname{index}\left(\Lambda^{1,0}, \text { Dolbeault }\right) \\
& =-c \int_{M} \operatorname{ch}\left(\Lambda^{1,0}\right)=c \int_{M} \operatorname{ch}\left(T^{\prime} M\right)=c \int_{M} e(M) .
\end{aligned}
$$

The Gauss-Bonnet theorem (or the normalization of Lemma 2.3.3) implies that the normalizing constant $c=1$ and gives another equivalent proof of the Riemann-Roch formula.

There are many applications of the Riemann-Roch theorem. We present a few in dimension 4 to illustrate some of the techniques involved. If $\operatorname{dim} M=4$, then we showed earlier that:

$$
c_{2}\left(T^{\prime} M\right)=e_{2}(T M) \quad \text { and } \quad\left(2 c_{2}-c_{1}^{2}\right)\left(T^{\prime} M\right)=p_{1}\left(T^{\prime} M\right)
$$

The Riemann-Roch formula expresses the arithmetic genus of $M$ in terms of $c_{2}$ and $c_{1}^{2}$. Consequently, there is a formula:

$$
\begin{aligned}
\operatorname{arithmetic} \text { genus } & =a_{1} \int_{M} e_{2}(T M)+a_{2} \int_{M} p_{1}(T M) / 3 \\
& =a_{1} \chi(M)+a_{2} \operatorname{signature}(M)
\end{aligned}
$$

where $a_{1}$ and $a_{2}$ are universal constants. If we consider the manifolds $S^{2} \times S^{2}$ and $\mathbf{C} P_{2}$ we derive the equations:

$$
1=a_{1} \cdot 4+a_{2} \cdot 0 \quad \text { and } \quad 1=a_{1} \cdot 3+a_{2} \cdot 1
$$

so that $a_{1}=a_{2}=1 / 4$ which proves:
Lemma 3.5.6. If $M$ is an almost complex manifold of real dimension 4, then:

$$
\operatorname{arithmetic} \operatorname{genus}(M)=\{\chi(M)+\operatorname{signature}(M)\} / 4
$$

Since the arithmetic genus is always an integer, we can use this result to obtain some non-integrability results:

Corollary 3.5.7. The following manifolds do not admit almost complex structures:
(a) $S^{4}$ (the four dimensional sphere).
(b) $\mathbf{C} P_{2}$ with the reversed orientation.
(c) $M_{1} \# M_{2}$ where the $M_{i}$ are 4-dimensional manifolds admitting almost complex structures (\# denotes connected sum).

Proof: We assume the contrary in each of these cases and attempt to compute the arithmetic genus:

$$
\begin{aligned}
\text { a.g. }\left(S^{4}\right) & =\frac{1}{4}(2+0)=\frac{1}{2} \\
\text { a.g. }\left(-\mathbf{C} P_{2}\right) & =\frac{1}{4}(3-1)=\frac{1}{2} \\
\text { a.g. }\left(M_{1} \# M_{2}\right) & =\frac{1}{4}\left(\chi\left(M_{1} \# d M_{2}\right)+\operatorname{sign}\left(M_{1} \# M_{2}\right)\right. \\
& =\frac{1}{4}\left(\chi\left(M_{1}\right)+\chi\left(M_{2}\right)-2+\operatorname{sign}\left(M_{1}\right)+\operatorname{sign}\left(M_{2}\right)\right) \\
& =\operatorname{a.g.}\left(M_{1}\right)+\operatorname{a.g} \cdot\left(M_{2}\right)-\frac{1}{2} .
\end{aligned}
$$

In none of these examples is the arithmetic genus an integer which shows the impossibility of constructing the desired almost complex structure.

### 3.6. A Review of Kaehler Geometry.

In the previous subsection, we proved the Riemann-Roch theorem using the $\operatorname{SPIN}_{c}$ complex. This was an essential step in the proof even if the manifold was holomorphic.

Let $M$ be holomorphic and let the coefficient bundle $V$ be holomorphic (i.e., the transition functions are holomorphic). We can define the Dolbeault complex directly by defining:

$$
\bar{\partial}_{V}: C^{\infty}\left(\Lambda^{0, q} \otimes V\right) \rightarrow C^{\infty}\left(\Lambda^{0, q+1} \otimes V\right) \quad \text { by } \bar{\partial}_{V}(\omega \otimes s)=\bar{\partial} \omega \otimes s
$$

where $s$ is a local holomorphic section to $V$. It is immediate $\bar{\partial}_{V} \bar{\partial}_{V}=0$ and that this is an elliptic complex; the index of this elliptic complex is just index ( $V$, Dolbeault) as defined previously.

We let $a_{n}(x, V$, Dolbeault $)$ be the invariant of the heat equation for this elliptic complex; we use the notation $a_{j}(x$, Dolbeault) when $V$ is the trivial bundle. Then:

Remark 3.6.1. Let $m=2 n>2$. Then there exists a unitary Riemannian metric on the m-torus (with its usual complex structure) and a point $x$ such that:
(a) $a_{j}(x$, Dolbeault $) \neq 0$ for $j$ even and $j \geq n$,
(b) $a_{m}(x$, Dolbeault $) \neq T d_{n}$ where both are viewed as scalar invariants.

The proof of this is quite long and combinatorial and is explained elsewhere; we simply present the result to demonstrate that it is not in general possible to prove the Riemann-Roch theorem directly by heat equation methods.

The difficulty is that the metric and the complex structure do not fit together properly. Choose local holomorphic coordinates $z=\left(z_{1}, \ldots, z_{n}\right)$ and extend the metric $G$ to be Hermitian on $T(M) \otimes \mathbf{C}$ so that $T^{\prime}$ and $T^{\prime \prime}$ are orthogonal. We define:

$$
g_{j \bar{k}}=G\left(\partial / \partial z_{j}, \partial / \partial z_{k}\right)
$$

then the matrix $g_{j \bar{k}}$ is a positive definite Hermitian matrix which determines the original metric on $T(M)$ :

$$
\begin{aligned}
d s^{2} & =2 \sum_{j, k} g_{j \bar{k}} d z^{j} \cdot d \bar{z}^{k} \\
G\left(\partial / \partial x_{j}, \partial / \partial x_{k}\right) & =G\left(\partial / \partial y_{j}, \partial / \partial y_{k}\right)=\left(g_{j \bar{k}}+g_{k \bar{\jmath}}\right) \\
G\left(\partial / \partial x_{j}, \partial / \partial y_{k}\right) & =-G\left(\partial / \partial y_{j}, \partial / \partial x_{k}\right)=\frac{1}{i}\left(g_{j \bar{k}}-g_{k \bar{\jmath}}\right)
\end{aligned}
$$

We use this tensor to define the Kaehler 2-form:

$$
\Omega=i \sum_{j, k} g_{j \bar{k}} d z^{j} \wedge d \bar{z}^{k}
$$

This is a real 2-form which is defined by the identity:

$$
\Omega[X, Y]=-G(X, J Y)
$$

(This is a slightly different sign convention from that sometimes followed.) The manifold $M$ is said to be Kaehler if $\Omega$ is closed and the heat equation gives a direct proof of the Riemann-Roch theorem for Kaehler manifolds.

We introduce variables:

$$
g_{j \bar{k} / l} \quad \text { and } \quad g_{j \bar{k} / \bar{l}}
$$

for the jets of the metric. On a Riemannian manifold, we can always find a coordinate system in which all the 1-jets of the metric vanish at a point; this concept generalizes as follows:

Lemma 3.6.2. Let $M$ be a holomorphic manifold and let $G$ be a unitary metric on $M$. The following statements are equivalent:
(a) The metric is Kaeher (i.e., $d \Omega=0$ ).
(b) For every $z_{0} \in M$ there is a holomorphic coordinate system $Z$ centred at $z_{0}$ so that $g_{j \bar{k}}(Z, G)\left(z_{0}\right)=\delta_{j k}$ and $g_{j \bar{k} / l}(Z, G)\left(z_{0}\right)=0$.
(c) For every $z_{0} \in M$ there is a holomorphic coordinate system $Z$ centred at $z_{0}$ so $g_{j \bar{k}}(Z, G)\left(z_{0}\right)=\delta_{j k}$ and so that all the 1 -jets of the metric vanish at $z_{0}$.

Proof: We suppose first the metric is Kaehler and compute:

$$
d \Omega=i \sum_{j, k}\left\{g_{j \bar{k} / l} d z^{l} \wedge d z^{j} \wedge d \bar{z}^{k}+g_{j \bar{k} / \bar{l}} d \bar{z}^{l} \wedge d z^{j} \wedge d \bar{z}^{k}\right\}
$$

Thus Kaehler is equivalent to the conditions:

$$
g_{j \bar{k} / l}-g_{l \bar{k} / j}=g_{j \bar{k} / \bar{l}}-g_{j \bar{l} / \bar{k}}=0
$$

By making a linear change of coordinates, we can assume that the holomorphic coordinate system is chosen to be orthogonal at the center $z_{0}$. We let

$$
z_{j}^{\prime}=z_{j}+\sum c_{k l}^{j} z_{k} z_{l} \quad\left(\text { where } c_{k l}^{j}=c_{l k}^{j}\right)
$$

and compute therefore:

$$
\begin{aligned}
d z_{j}^{\prime} & =d z_{j}+2 \sum c_{k l}^{j} z_{l} d z_{k} \\
\partial / \partial z_{j}^{\prime} & =\partial / \partial z_{j}-2 \sum c_{j l}^{k} z_{l} \partial / \partial z_{k}+O\left(z^{2}\right) \\
g_{j \bar{k}}^{\prime} & =g_{j \bar{k}}-2 \sum c_{j l}^{k} z_{l}+\text { terms in } \bar{z}+O\left(|z|^{2}\right) \\
g_{j \bar{k} / l}^{\prime} & =g_{j \bar{k} / l}-2 c_{j l}^{k} \quad \text { at } z_{0}
\end{aligned}
$$

We define $c_{j l}^{k}=\frac{1}{2} g_{j \bar{k} / l}$ and observe that the Kaehler condition shows this is symmetric in the indices $(j, l)$ so (a) implies (b). To show (b) implies (c) we simply note that $g_{j \bar{k} / \bar{l}}=\overline{g_{k \bar{\jmath} / l}}=0$ at $z_{0}$. To prove (c) implies (a) we observe $d \Omega\left(z_{0}\right)=0$ so since $z_{0}$ was arbitrary, $d \Omega \equiv 0$.

If $M$ is Kaehler, then $\delta \Omega$ is linear in the 1 -jets of the metric when we compute with respect to a holomorphic coordinate system. This implies that $\delta \Omega=0$ so $\Omega$ is harmonic. We compute that $\Omega^{k}$ is also harmonic for $1 \leq k \leq n$ and that $\Omega^{n}=c \cdot$ dvol is a multiple of the volume form (and in particular is non-zero). This implies that if $x$ is the element in $H^{2}(M ; \mathbf{C})$ defined by $\Omega$ using the Hodge decomposition theorem, then $1, x, \ldots, x^{n}$ all represent non-zero elements in the cohomology ring of $M$. This can be used to show that there are topological obstructions to constructing Kaehler metrics:

REMARK 3.6.3. If $M=S^{1} \times S^{m-1}$ for $m$ even and $m \geq 4$ then this admits a holomorphic structure. No holomorphic structure on $M$ admits a Kaehler metric.

Proof: Since $H^{2}(M ; \mathbf{C})=0$ it is clear $M$ cannot admit a Kaehler metric. We construct holomorphic structures on $M$ as follows: let $\lambda \in \mathbf{C}$ with $|\lambda|>1$ and let $M_{\lambda}=\left\{\mathbf{C}^{n}-0\right\} / \lambda$ (where we identify $z$ and $w$ if $z=\lambda^{k} w$ for some $k \in \mathbf{Z}$ ). The $M_{\lambda}$ are all topologically $S^{1} \times S^{m-1}$; for example if $\lambda \in \mathbf{R}$ and we introduce spherical coordinates $(r, \theta)$ on $\mathbf{C}^{n}=\mathbf{R}^{m}$ then we are identifying $(r, \theta)$ and $\left(\lambda^{k} r, \theta\right)$. If we let $t=\log r$, then we are identifying $t$ with $t+k \log \lambda$ so the manifold is just $[1, \log \lambda] \times S^{m-1}$ where we identify the endpoints of the interval. The topological identification of $M_{\lambda}$ for other $\lambda$ is similar.

We now turn to the problem of constructing a Kaehler metric on $\mathbf{C} P_{n}$. Let $L$ be the tautological line bundle over $\mathbf{C} P_{n}$ and let $x=-c_{1}(L) \in$ $\Lambda^{1,1}\left(\mathbf{C} P_{n}\right)$ be the generator of $H^{2}\left(\mathbf{C} P_{n} ; \mathbf{Z}\right)$ discussed earlier. We expand $x$ in local coordinates in the form:

$$
x=\frac{i}{2 \pi} \sum_{j, k} g_{j \bar{k}} d z^{j} \wedge d \bar{z}^{k}
$$

and define $G\left(\partial / \partial z_{k}, \partial / \partial z_{k}\right)=g_{j \bar{k}}$. This gives an invariantly defined form called the Fubini-Study metric on $T_{c}\left(\mathbf{C} P_{n}\right)$. We will show $G$ is positive definite and defines a unitary metric on the real tangent space. Since the Kaehler form of $G, \Omega=2 \pi x$, is a multiple of $x$, we conclude $d \Omega=0$ so $G$ will be a Kaehler metric.

The 2-form $x$ is invariant under the action of $\mathrm{U}(n+1)$. Since $\mathrm{U}(n+1)$ acts transitively on $\mathbf{C} P_{n}$, it suffices to show $G$ is positive definite and symmetric at a single point. Using the notation of section 2.3, let $\mathrm{U}_{n}=\left\{z \in \mathbf{C} P_{n}\right.$ : $\left.z_{n+1}(z) \neq 0\right\}$. We identify $\mathrm{U}_{n}$ with $\mathbf{C}^{n}$ by identifying $z \in \mathbf{C}^{n}$ with the line
in $\mathbf{C}^{n+1}$ through the point $(z, 1)$. Then:

$$
x=\frac{i}{2 \pi} \partial \bar{\partial} \log \left(1+|z|^{2}\right)
$$

so at the center $z=0$ we compute that:

$$
x=\frac{i}{2 \pi} \sum_{j} d z^{j} \wedge d \bar{z}^{j} \quad \text { so that } \quad g_{j \bar{k}}=\delta_{j \bar{k}}
$$

If $G^{\prime}$ is any other unitary metric on $\mathbf{C} P_{n}$ which is invariant under the action of $\mathrm{U}(n+1)$, then $G^{\prime}$ is determined by its value on $T_{c}\left(\mathbf{C} P_{n}\right)$ at a single point. Since $\mathrm{U}(n)$ preserves the origin of $\mathrm{U}_{n}=\mathbf{C}^{n}$, we conclude $G^{\prime}=c G$ for some $c>0$. This proves:

Lemma 3.6.4.
(a) There exists a unitary metric $G$ on $\mathbf{C} P_{n}$ which is invariant under the action of $\mathrm{U}(n+1)$ such that the Kaehler form of $G$ is given by $\Omega=2 \pi x$ where $x=-c_{1}(L)\left(L\right.$ is the tautological bundle over $\left.\mathbf{C} P_{n}\right) . G$ is a Kaehler metric.
(b) $\int_{\mathbf{C} P_{n}} x^{n}=1$.
(c) If $G^{n}$ is any other unitary metric on $\mathbf{C} P_{n}$ which is invariant under the action of $\mathrm{U}(n+1)$ then $G^{\prime}=c G$ for some constant $c>0$.

A holomorphic manifold $M$ is said to be Hodge if it admits a Kaehler metric $G$ such that the Kaehler form $\Omega$ has the form $\Omega=c_{1}(L)$ for some holomorphic line bundle $L$ over $M$. Lemma 3.6 .4 shows $\mathbf{C} P_{n}$ is a Hodge manifold. Any submanifold of a Hodge manifold is again Hodge where the metric is just the restriction of the given Hodge metric. Thus every algebraic variety is Hodge. The somewhat amazing fact is that the converse is true:

Remark 3.6.5. A holomorphic manifold $M$ is an algebraic variety (i.e., is holomorphically equivalent to a manifold defined by algebraic equations in $\mathbf{C} P_{n}$ for some $n$ ) if and only if it admits a Hodge metric.

We shall not prove this remark but simply include it for the sake of completeness. We also note in passing that there are many Kaehler manifolds which are not Hodge. The Riemann period relations give obstructions on complex tori to those tori being algebraic.

We now return to our study of Kaehler geometry. We wish to relate the $\bar{\partial}$ cohomology of $M$ to the ordinary cohomology. Define:

$$
H^{p, q}(M)=\operatorname{ker} \bar{\partial} / \operatorname{im} \bar{\partial} \quad \text { in bi-degree }(p, q)
$$

Since the Dolbeault complex is elliptic, these groups are finite dimensional using the Hodge decomposition theorem.

Lemma 3.6.6. Let $M$ be Kaehler and let $\Omega=i \sum g_{j \bar{k}} d z^{j} \wedge d \bar{z}^{k}$ be the Kaehler form. Let $d=\partial+\bar{\partial}$ and let $\delta=\delta^{\prime}+\delta^{\prime \prime}$. Then we have the identities:
(a) $\bar{\partial} \operatorname{int}(\Omega)-\operatorname{int}(\Omega) \bar{\partial}=i \delta^{\prime}$.
(b) $\bar{\partial} \delta^{\prime}+\delta^{\prime} \bar{\partial}=\partial \delta^{\prime \prime}+\delta^{\prime \prime} \partial=0$.
(c) $d \delta+\delta d=2\left(\partial \delta^{\prime}+\delta^{\prime} \partial\right)=2\left(\bar{\partial} \delta^{\prime \prime}+\delta^{\prime \prime} \bar{\partial}\right)$.
(d) We can decompose the de Rham cohomology in terms of the Dolbeault cohomology:

$$
H^{n}(M ; \mathbf{C})=\bigoplus_{p+q=n} H^{p, q}(M ; \mathbf{C})
$$

(e) There are isomorphisms $H^{p, q} \simeq H^{q, p} \simeq H^{n-p, n-q} \simeq H^{n-q, n-p}$.

Proof: We first prove (a). We define:

$$
A=\bar{\partial} \operatorname{int}(\Omega)-\operatorname{int}(\Omega) \bar{\partial}-i \delta^{\prime}
$$

If we can show that $A$ is a $0^{\text {th }}$ order operator, then $A$ will be a functorially defined endomorphism linear in the 1-jets of the metric. This implies $A=0$ by Lemma 3.6.2. Thus we must check $\sigma_{L}(A)=0$. We choose an orthonormal basis $\left\{e_{j}\right\}$ for $T(M)=T^{*}(M)$ so $J e_{j}=e_{j+n}, 1 \leq j \leq n$. By Lemma 3.5.2, it suffices to show that:

$$
\operatorname{ext}\left(e_{j}-i e_{j+n}\right) \operatorname{int}(\Omega)-\operatorname{int}(\Omega) \operatorname{ext}\left(e_{j}-i e_{j+n}\right)=-i \cdot \operatorname{int}\left(e_{j}-i e_{j+n}\right)
$$

for $1 \leq j \leq n$. Since $g_{j \bar{k}}=\frac{1}{2} \delta_{j \bar{k}}, \Omega$ is given by

$$
\Omega=\frac{i}{2} \sum_{j}\left(e_{j}+i e_{j+n}\right) \wedge\left(e_{j}-i e_{j+n}\right)=\sum_{j} e_{j} \wedge e_{j+n} .
$$

The commutation relations:

$$
\operatorname{int}\left(e_{k}\right) \operatorname{ext}\left(e_{l}\right)+\operatorname{ext}\left(e_{l}\right) \operatorname{int}\left(e_{k}\right)=\delta_{k, l}
$$

imply that $\operatorname{int}\left(e_{k} \wedge e_{k+n}\right)=-\operatorname{int}\left(e_{k}\right) \operatorname{int}\left(e_{k+n}\right)$ commutes with $\operatorname{ext}\left(e_{j}-\right.$ $\left.i e_{j+n}\right)$ for $k \neq j$ so these terms disappear from the commutator. We must show:

$$
\begin{aligned}
\operatorname{ext}\left(e_{j}-i e_{j+n}\right) \operatorname{int}\left(e_{j} \wedge e_{j+n}\right)-\operatorname{int}\left(e_{j} \wedge e_{j+n}\right) & \operatorname{ext}\left(e_{j}-i e_{j+n}\right) \\
& =-i \cdot \operatorname{int}\left(e_{j}-i e_{j+n}\right)
\end{aligned}
$$

This is an immediate consequence of the previous commutation relations for int and ext together with the identity $\operatorname{int}\left(e_{j}\right) \operatorname{int}\left(e_{k}\right)+\operatorname{int}\left(e_{k}\right) \operatorname{int}\left(e_{j}\right)=0$. This proves (a).

From (a) we compute:

$$
\bar{\partial} \delta^{\prime}+\delta^{\prime} \bar{\partial}=i(-\bar{\partial} \operatorname{int}(\Omega) \bar{\partial})+i(\bar{\partial} \operatorname{int}(\Omega) \bar{\partial})=0
$$

and by taking complex conjugates $\partial \delta^{\prime \prime}+\delta^{\prime \prime} \partial=0$. This shows (b). Thus
$\Delta=(d+\delta)^{2}=d \delta+\delta d=\left\{\partial \delta^{\prime}+\delta^{\prime} \partial\right\}+\left\{\bar{\partial} \delta^{\prime \prime}+\delta^{\prime \prime} \partial\right\}=\left\{\partial+\delta^{\prime}\right\}^{2}+\left\{\bar{\partial}+\delta^{\prime \prime}\right\}^{2}$
since all the cross terms cancel. We apply (a) again to compute:

$$
\begin{aligned}
\partial \bar{\partial} \operatorname{int}(\Omega)-\partial \operatorname{int}(\Omega) \bar{\partial} & =i \partial \delta^{\prime} \\
\bar{\partial} \operatorname{int}(\Omega) \partial-\operatorname{int}(\Omega) \bar{\partial} \partial & =i \delta^{\prime} \partial
\end{aligned}
$$

which yields the identity:

$$
i\left(\partial \delta^{\prime}+\delta^{\prime} \partial\right)=\partial \bar{\partial} \operatorname{int}(\Omega)-\partial \operatorname{int}(\Omega) \bar{\partial}+\bar{\partial} \operatorname{int}(\Omega) \partial-\operatorname{int}(\Omega) \bar{\partial} \partial
$$

Since $\Omega$ is real, we take complex conjugate to conclude:

$$
\begin{aligned}
-i\left(\bar{\partial} \delta^{\prime \prime}+\delta^{\prime \prime} \bar{\partial}\right) & =\bar{\partial} \partial \operatorname{int}(\Omega)-\bar{\partial} \operatorname{int}(\Omega) \partial+\partial \operatorname{int}(\Omega) \bar{\partial}-\operatorname{int}(\Omega) \partial \bar{\partial} \\
& =-i\left(\partial \delta^{\prime}+\delta^{\prime} \partial\right)
\end{aligned}
$$

This shows $\left(\bar{\partial} \delta^{\prime \prime}+\delta^{\prime \prime} \bar{\partial}\right)=\left(\partial \delta^{\prime}+\delta^{\prime} \partial\right)$ so $\Delta=2\left(\bar{\partial} \delta^{\prime \prime}+\delta^{\prime \prime} \bar{\partial}\right)$ which proves (c).
$H^{n}(M ; \mathbf{C})$ denotes the de Rham cohomology groups of $M$. The Hodge decomposition theorem identifies these groups with the null-space of $(d+\delta)$. Since $\mathrm{N}(d+\delta)=\mathrm{N}\left(\bar{\partial}+\delta^{\prime \prime}\right)$, this proves (d). Finally, taking complex conjugate and applying $*$ induces the isomorphisms:

$$
H^{p, q} \simeq H^{q, p} \simeq H^{n-q, n-p} \simeq H^{n-p, n-q} .
$$

In particular $\operatorname{dim} H^{1}(M ; \mathbf{C})$ is even if $M$ admits a Kaehler metric. This gives another proof that $S^{1} \times S^{2 n-1}$ does not admit a Kaehler metric for $n>1$.

The duality operations of (e) can be extended to the Dolbeault complex with coefficients in a holomorphic bundle $V$. If $V$ is holomorphic, the transition functions are holomorphic and hence commute with $\bar{\partial}$. We define:

$$
\bar{\partial}: C^{\infty}\left(\Lambda^{p, q} \otimes V\right) \rightarrow\left(\Lambda^{p, q+1} \otimes V\right)
$$

by defining $\bar{\partial}(\omega \otimes s)=\bar{\partial} \omega \otimes s$ relative to any local holomorphic frame for $V$. This complex is equivalent to the one defined in section 3.5 and defines cohomology classes:

$$
H^{p, q}(V)=\operatorname{ker} \bar{\partial} / \operatorname{im} \bar{\partial} \quad \text { on } C^{\infty}\left(\Lambda^{p, q} \otimes V\right)
$$

It is immediate from the definition that:

$$
H^{p, q}(V)=H^{0, q}\left(\Lambda^{p, 0} \otimes V\right)
$$

We define

$$
\operatorname{index}(V, \bar{\partial})=\operatorname{index}(V, \text { Dolbeault })=\sum(-1)^{q} \operatorname{dim} H^{0, q}(V)
$$

to be the index of this elliptic complex.
There is a natural duality (called Serre duality) which is conjugate linear $\bar{*}: H^{p, q}(V) \rightarrow H^{n-p, n-q}\left(V^{*}\right)$, where $V^{*}$ is the dual bundle. We put a unitary fiber metric on $V$ and define:

$$
T\left(s_{1}\right)\left(s_{2}\right)=s_{2} \cdot s_{1}
$$

so $T: V \rightarrow V^{*}$ is conjugate linear. If $*: \Lambda^{k} \rightarrow \Lambda^{m-k}$ is the ordinary Hodge operator, then it extends to define $*: \Lambda^{p, q} \rightarrow \Lambda^{n-q, n-p}$. The complex conjugate operator defines $\bar{*}: \Lambda^{p, q} \rightarrow \Lambda^{n-p, n-q}$. We extend this to have coefficients in $V$ by defining:

$$
\bar{*}(\omega \otimes s)=\bar{\star} \omega \otimes T s
$$

so $\bar{*}: \Lambda^{p, q} \otimes V \rightarrow \Lambda^{n-p, n-q} \otimes V^{*}$ is conjugate linear. Then:
Lemma3.6.7. Let $\delta_{V}^{\prime \prime}$ be the adjoint of $\bar{\partial}_{V}$ on $C^{\infty}\left(\Lambda^{p, *} \otimes V\right)$ with respect to a unitary metric on $M$ and a Hermitian metric on $V$. Then
(a) $\delta^{\prime \prime}=-\bar{*} \bar{\partial}_{V} \bar{*}$.
(b) $\begin{gathered}\text { induces a conjugate linear isomorphism between the groups } H^{p, q}(V)\end{gathered}$ and $H^{n-p, n-q}\left(V^{*}\right)$.

Proof: It is important to note that this lemma, unlike the previous one, does not depend upon having a Kaehler metric. We suppose first $V$ is holomorphically trivial and the metric on $V$ is flat. Then $\Lambda \otimes V=\Lambda \otimes 1^{k}$ for some $k$ and we may suppose $k=1$ for the sake of simplicity. We noted

$$
\delta=-* d *=-\bar{*} d \bar{*}
$$

in section 1.5. This decomposes

$$
\delta=\delta^{\prime}+\delta^{\prime \prime}=-\bar{\star} \partial \bar{\not}-\bar{\star} \bar{\partial} \bar{\not} .
$$

Since $\bar{*}: \Lambda^{p, q} \rightarrow \Lambda^{n-p, n-q}$, we conclude $\delta^{\prime}=-\bar{\not} \partial \bar{\not}$ and $\delta^{\prime \prime}=-\bar{*} \bar{\partial} \bar{\not}$. This completes the proof if $V$ is trivial. More generally, we define

$$
A=\delta_{V}^{\prime \prime}+\bar{*} \bar{\partial} \bar{\not} .
$$

This must be linear in the 1 -jets of the metric on $V$. We can always normalize the choice of local holomorphic frame so the 1-jets vanish at a basepoint $z_{0}$ and thus $A=0$ in general which proves (a). We identify $H^{p, q}(V)$ with $\mathrm{N}\left(\bar{\partial}_{V}\right) \cap \mathrm{N}\left(\delta_{V}^{\prime \prime}\right)$ using the Hodge decomposition theorem. Thus $\bar{*}: H^{p, q}(V) \rightarrow H^{n-p, n-q}\left(V^{*}\right)$. Since $\bar{*}^{2}= \pm 1$, this completes the proof.

We defined $x=-c_{1}(L) \in \Lambda^{1,1}\left(\mathbf{C} P_{n}\right)$ for the standard metric on the tautological line bundle over $\mathbf{C} P_{n}$. In section 2.3 we showed that $x$ was harmonic and generates the cohomology ring of $\mathbf{C} P_{n}$. Since $x$ defines a Kaehler metric, called the Fubini-Study metric, Lemma 3.6.6 lets us decompose $H^{n}\left(\mathbf{C} P_{n} ; \mathbf{C}\right)=\bigoplus_{p+q=n} H^{p, q}\left(\mathbf{C} P_{n}\right)$. Since $x^{k} \in H^{k, k}$ this shows: Lemma 3.6.8. Let $\mathbf{C} P_{n}$ be given the Fubini-Study metric. Let $x=$ $-c_{1}(L)$. Then $H^{k, k}(M) \simeq \mathbf{C}$ is generated by $x^{k}$ for $1 \leq k \leq n . H^{p, q}(M)=$ 0 for $p \neq q$. Thus in particular,

$$
\operatorname{index}(\bar{\partial})=\text { arithmetic genus }=\sum(-1)^{k} \operatorname{dim} H^{0, k}=1 \quad \text { for } \mathbf{C} P_{n}
$$

Since the arithmetic genus is multiplicative with respect to products,

$$
\operatorname{index}(\bar{\partial})=1 \quad \text { for } \mathbf{C} P_{n_{1}} \times \cdots \times \mathbf{C} P_{n_{k}}
$$

We used this fact in the previous subsection to derive the normalizing constants in the Riemann-Roch formula.

We can now present another proof of the Riemann-Roch theorem for Kaehler manifolds which is based directly on the Dolbeault complex and not on the $\mathrm{SPIN}_{c}$ complex. Let $Z$ be a holomorphic coordinate system and let

$$
g_{u \bar{v}}=G\left(\partial / \partial z_{u}, \partial / \partial z_{v}\right)
$$

represent the components of the metric tensor. We introduce additional variables:

$$
\begin{aligned}
& g_{u \bar{v} / \alpha \bar{\beta}}=d_{z}^{\alpha} d_{\bar{z}}^{\beta} g_{u \bar{v}} \\
& \quad \quad \text { for } 1 \leq u, v \leq n, \alpha=\left(\alpha_{1}, \ldots, \alpha_{n}\right), \text { and } \beta=\left(\beta_{1}, \ldots, \beta_{n}\right)
\end{aligned}
$$

These variables represent the formal derivatives of the metric.
We must also consider the Dolbeault complex with coefficients in a holomorphic vector bundle $V$. We choose a local fiber metric $H$ for $V$. If $s=\left(s_{1}, \ldots, s_{k}\right)$ is a local holomorphic frame for $V$, we introduce variables:

$$
h_{p, \bar{q}}=H\left(s_{p}, s_{q}\right), \quad h_{p, \bar{q} / \alpha \bar{\beta}}=d_{z}^{\alpha} d_{\bar{z}}^{\beta} h_{p \bar{q}} \quad \text { for } 1 \leq p, q \leq k=\operatorname{dim} V
$$

If $Z$ is a holomorphic coordinate system centered at $z_{0}$ and if $s$ is a local holomorphic frame, we normalize the choice so that:

$$
g_{u \bar{v}}(Z, G)\left(z_{0}\right)=\delta_{u, v} \quad \text { and } \quad h_{p \bar{q}}(s, H)\left(z_{0}\right)=\delta_{p, q}
$$

The coordinate frame $\left\{\partial / \partial z_{u}\right\}$ and the holomorphic frame $\left\{s_{p}\right\}$ for $V$ are orthogonal at $z_{0}$. Let $\mathcal{R}$ be the polynomial algebra in these variables. If $P \in \mathcal{R}$ we can evaluate $P(Z, G, s, H)\left(z_{0}\right)$ once a holomorphic frame $Z$ is chosen and a holomorphic frame $s$ is chosen. We say $P$ is invariant if $P(Z, G, s, H)\left(z_{0}\right)=P(G, H)\left(z_{0}\right)$ is independent of the particular $(Z, s)$ chosen; let $\mathcal{R}_{n, k}^{c}$ denote the sub-algebra of all invariant polynomials. As before, we define:

$$
\operatorname{ord}\left(g_{u, \bar{v} / \alpha \bar{\beta}}\right)=\operatorname{ord}\left(s_{p \bar{q} / \alpha \bar{\beta}}\right)=|\alpha|+|\beta| .
$$

We need only consider those variables of positive order as $g_{u, \bar{v}}=\delta_{j \bar{v}}$ and $s_{p, \bar{q}}=\delta_{p, q}$ at $z_{0}$. If $\mathcal{R}_{n, \nu, k}^{c}$ denotes the subspace of invariant polynomials homogenous of order $\nu$ in the jets of the metrics $(G, H)$, then there is a direct sum decomposition:

$$
\mathcal{R}_{n, k}^{c}=\bigoplus \mathcal{R}_{n, \nu, k}^{c}
$$

In the real case in studying the Euler form, we considered the restriction map from manifolds of dimension $m$ to manifolds of dimension $m-1$. In the complex case, there is a natural restriction map

$$
r: \mathcal{R}_{n, \nu, k}^{c} \rightarrow \mathcal{R}_{n-1, \nu, k}^{c}
$$

which lowers the complex dimension by one (and the real dimension by two). Algebraically, $r$ is defined as follows: let

$$
\begin{aligned}
\operatorname{deg}_{k}\left(g_{u \bar{v} / \alpha \bar{\beta}}\right) & =\delta_{k, u}+\alpha(k) & & \operatorname{deg}_{\bar{k}}\left(g_{u \bar{v} / \alpha \bar{\beta}}\right)=\delta_{k, v}+\beta(k) \\
\operatorname{deg}_{k}\left(h_{p \bar{q} / \alpha \bar{\beta}}\right) & =\alpha(k) & & \operatorname{deg}_{\bar{k}}\left(h_{p \bar{q} / \alpha \bar{\beta}}\right)=\beta(k) .
\end{aligned}
$$

We define:

$$
\begin{aligned}
& r\left(g_{u \bar{v} / \alpha \bar{\beta}}\right)= \begin{cases}g_{u \bar{v} / \alpha \bar{\beta}} & \text { if } \operatorname{deg}_{n}\left(g_{u \bar{v} / \alpha \bar{\beta}}\right)+\operatorname{deg}_{\bar{n}}\left(g_{u \bar{v} / \alpha \bar{\beta}}\right)=0 \\
0 & \text { if } \operatorname{deg}_{n}\left(g_{u \bar{v} / \alpha \bar{\beta}}\right)+\operatorname{deg}_{\bar{n}}\left(g_{u \bar{v} / \alpha \bar{\beta}}\right)>0\end{cases} \\
& r\left(h_{p \bar{q} / \alpha \bar{\beta}}\right)= \begin{cases}h_{p \bar{q} / \alpha \bar{\beta}} & \text { if } \operatorname{deg}_{n}\left(h_{p \bar{q} / \alpha \bar{\beta}}\right)+\operatorname{deg}_{\bar{n}}\left(h_{p \bar{q} / \alpha \bar{\beta}}\right)=0 \\
0 & \text { if } \operatorname{deg}_{n}\left(h_{p \bar{q} / \alpha \bar{\beta}}\right)+\operatorname{deg}_{\bar{n}}\left(h_{p \bar{q} / \alpha \bar{\beta}}\right)>0 .\end{cases}
\end{aligned}
$$

This defines a map $r: \mathcal{R}_{n} \rightarrow \mathcal{R}_{n-1}$ which is an algebra morphism; we simply set to zero those variables which do not belong to $\mathcal{R}_{n-1}$. It is immediate that $r$ preserves both invariance and the order of a polynomial. Geometrically, we are considering manifolds of the form $M^{n}=M^{n-1} \times T_{2}$ where $T_{2}$ is the flat two torus, just as in the real case we considered manifolds $M^{m}=M^{m-1} \times S^{1}$.

Let $\mathcal{R}_{n, p, k}^{c h}$ be the space of $p$-forms generated by the Chern forms of $T_{c}(M)$ and by the Chern forms of $V$. We take the holomorphic connection defined by the metrics $G$ and $H$ on $T_{c}(M)$ and on $V$. If $P \in \mathcal{R}_{n, m, k}^{c h}$ then $* P$ is a scalar invariant. Since $P$ vanishes on product metrics of the form $M^{n-1} \times T_{2}$ which are flat in one holomorphic direction, $r(* P)=0$. This is the axiomatic characterization of the Chern forms which we shall need.

Theorem 3.6.9. Let the complex dimension of $M$ be $n$ and let the fiber dimension of $V$ be $k$. Let $P \in \mathcal{R}_{n, \nu, k}^{c}$ and suppose $r(P)=0$. Then:
(a) If $\nu<2 n$ then $P(G, H)=0$ for all Kaehler metrics $G$.
(b) If $\nu=2 n$ then there exists a unique $Q \in \mathcal{R}_{n, m, k}^{c h}$ so $P(G, H)=$ * $Q(G, H)$ for all Kaehler metrics $G$.

As the proof of this theorem is somewhat technical, we shall postpone the proof until section 3.7. This theorem is false if we don't restrict to Kaehler metrics. There is a suitable generalization to form valued invariants. We can apply Theorem 3.6.9 to give a second proof of the Riemann-Roch theorem:

Theorem 3.6.10. Let $V$ be a holomorphic bundle over $M$ with fiber metric $H$ and let $G$ be a Kaehler metric on $M$. Let $\bar{\partial}_{V}: C^{\infty}\left(\Lambda^{0, q} \otimes V\right) \rightarrow$ $C^{\infty}\left(\Lambda^{0, q+1} \otimes V\right)$ denote the Dolbeault complex with coefficients in $V$ and let $a_{\nu}\left(x, \bar{\partial}_{V}\right)$ be the invariants of the heat equation. $a_{\nu} \in \mathcal{R}_{n, \nu, k}^{c}$ and

$$
\int_{M} a_{\nu}\left(x, \bar{\partial}_{V}\right) \operatorname{dvol}(x)= \begin{cases}0 & \text { if } \nu \neq 2 n \\ \operatorname{index}\left(\bar{\partial}_{V}\right) & \text { if } \nu=2 n\end{cases}
$$

Then
(a) $a_{\nu}\left(x, \bar{\partial}_{V}\right)=0$ for $\nu<2 n$.
(b) $a_{2 n}\left(x, \bar{\partial}_{V}\right)=*\left\{T d\left(T_{c} M\right) \wedge c h(V)\right\}_{m}$.

Proof: We note this implies the Riemann-Roch formula. By remark 3.6.1, we note this theorem is false in general if the metric is not assumed to be Kaehler. The first assertions of the theorem (including the homogeneity) follow from the results of Chapter 1 , so it suffices to prove (a) and (b). The Dolbeault complex is multiplicative under products $M=M^{n-1} \times T_{2}$. If we take the product metric and assume $V$ is the pull-back of a bundle over $M_{n-1}$ then the natural decomposition:

$$
\Lambda^{0, q} \otimes V \simeq \Lambda^{0, q}\left(M^{n-1}\right) \otimes V \oplus \Lambda^{0, q-1}\left(M^{n-1}\right) \otimes V
$$

shows that $r\left(a_{\nu}\right)=0$. By Theorem 3.6.9, this shows $a_{\nu}=0$ for $\nu<2 n$ proving (a).

In the limiting case $\nu=2 n$, we conclude $* a_{2 n}$ is a characteristic $2 n$ form. We first suppose $m=1$. Any one dimensional complex manifold is Kaehler. Then

$$
\operatorname{index}\left(\bar{\partial}_{V}\right)=a_{1} \int_{M} c_{1}\left(T_{c}(M)\right)+a_{2} \int_{M} c_{1}(V)
$$

where $a_{1}$ and $a_{2}$ are universal constants to be determined. Lemma 3.6.6 implies that $\operatorname{dim} H^{0,0}=\operatorname{dim} H^{0}=1=\operatorname{dim} H^{2}=\operatorname{dim} H^{1,1}$ while $\operatorname{dim} H^{1,0}=$ $\operatorname{dim} H^{0,1}=g$ where $g$ is the genus of the manifold. Consequently:

$$
\chi(M)=2-2 g=2 \operatorname{index}(\bar{\partial}) .
$$

We specialize to the case $M=S^{2}$ and $V=1$ is the trivial bundle. Then $g=0$ and

$$
1=\operatorname{index}\left(\bar{\partial}_{V}\right)=a_{1} \int_{M} c_{1}\left(T_{c}(M)\right)=2 a_{1}
$$

so $a_{1}=\frac{1}{2}$. We now take the line bundle $V=\Lambda^{1,0}$ so

$$
\operatorname{index}\left(\bar{\partial}_{V}\right)=g-1=-1=1+a_{2} \int_{M} c_{1}\left(T_{c}^{*}\right)=1-2 a_{2}
$$

so that $a_{2}=1$. Thus the Riemann-Roch formula in dimension 1 becomes:

$$
\operatorname{index}\left(\bar{\partial}_{V}\right)=\frac{1}{2} \int_{M} c_{1}\left(T_{c}(M)\right)+\int_{M} c_{1}(V)
$$

We now use the additive nature of the Dolbeault complex with respect to $V$ and the multiplicative nature with respect to products $M_{1} \times M_{2}$ to conclude that the characteristic $m$-form $* a_{m}$ must have the form:

$$
\left\{T_{d}^{\prime}\left(T_{c} M\right) \wedge \operatorname{ch}(V)\right\}
$$

where we use the fact the normalizing constant for $c_{1}(V)$ is 1 if $m=1$. We now use Lemma 2.3.5(a) together with the fact that index $\left(\bar{\partial}_{V}\right)=1$ for products of complex projective spaces to show $T d^{\prime}=T d$. This completes the proof.

### 3.7. An Axiomatic Characterization <br> Of the Characteristic Forms for <br> Holomorphic Manifolds with Kaehler Metrics.

In subsection 3.6 we gave a proof of the Riemann-Roch formula for Kaehler manifolds based on Theorem 3.6 .9 which gave an axiomatic characterization of the Chern forms. This subsection will be devoted to giving the proof of Theorem 3.6.9. The proof is somewhat long and technical so we break it up into a number of steps to describe the various ideas which are involved.

We introduce the notation:

$$
g_{u_{0} v_{0} / u_{1} \ldots u_{j} \bar{v}_{1} \ldots \bar{v}_{k}}
$$

for the jets of the metric on $M$. Indices $(u, v, w)$ will refer to $T_{c}(M)$ and will run from 1 thru $n=\frac{m}{2}$. We use ( $\bar{u}, \bar{v}, \bar{w}$ ) for anti-holomorphic indices. The symbol "*" will refer to indices which are not of interest in some particular argument. We let $A_{0}$ denote a generic monomial. The Kaehler condition is simply the identity:

$$
g_{u \bar{v} / w}=g_{w \bar{v} / u} \quad \text { and } \quad g_{u \bar{v} / \bar{w}}=g_{u \bar{w} / \bar{v}}
$$

We introduce new variables:

$$
g\left(u_{0}, \ldots, u_{j} ; \bar{v}_{0}, \ldots, \bar{v}_{k}\right)=g_{u_{0} v_{0} / u_{1} \ldots u_{j} \bar{v}_{1} \ldots \bar{v}_{k}} .
$$

If we differentiate the Kaehler identity, then we conclude $g(\vec{u} ; \vec{v})$ is symmetric in $\vec{u}=\left(u_{0}, \ldots, u_{j}\right)$ and $\vec{v}=\left(v_{0}, \ldots, v_{j}\right)$. Consequently, we may also use the multi-index notation $g(\alpha ; \bar{\beta})$ to denote these variables. We define:
$\operatorname{ord}(g(\alpha ; \bar{\beta}))=|\alpha|+|\beta|-2, \quad \operatorname{deg}_{u} g(\alpha ; \bar{\beta})=\alpha(u), \quad \operatorname{deg}_{\bar{u}} g(\alpha ; \bar{\beta})=\beta(u)$.
We already noted in Lemma 3.6.2 that it was possible to normalize the coordinates so $g_{u \bar{v} / w}\left(z_{0}\right)=g_{u \bar{v} / \bar{w}}\left(z_{0}\right)=0$. The next lemma will permit us to normalize coordinates to arbitrarily high order modulo the action of the unitary group:

Lemma 3.7.1. Let $G$ be a Kaehler metric and let $z_{0} \in M$. Let $\nu \geq 2$ be given. Then there exists a holomorphic coordinate system $Z$ centered at $z_{0}$ so that
(a) $g_{u \bar{v}}(Z, G)\left(z_{0}\right)=\delta_{u v}$.
(b) $g_{u \bar{v} / \alpha}(Z, G)\left(z_{0}\right)=g_{u \bar{v} / \bar{\alpha}}(Z, G)\left(z_{0}\right)=0$ for $1 \leq|\alpha| \leq \nu-1$.
(c) $Z$ is unique modulo the action of the unitary group $\mathrm{U}(n)$ and modulo coordinate transformations of order $\nu+1$ in $z$.

Proof: We proceed by induction. It is clear $\mathrm{U}(n)$ preserves such coordinate systems. The case $\nu=2$ is just Lemma 3.6.2 and the uniqueness is clear. We now consider $W$ given which satisfies (a) and (b) for
$1 \leq|\alpha| \leq \nu-2$ and define a new coordinate system $Z$ by setting:

$$
w_{u}=z_{u}+\sum_{|\alpha|=\nu} C_{u, \alpha} z^{\alpha}
$$

where the constants $C_{u, \alpha}$ remain to be determined. Since this is the identity transformation up to order $\nu$, the $\nu-2$ jets of the metric are unchanged so conditions (a) and (b) are preserved for $|\beta|<\nu-1$. We compute:

$$
\frac{\partial}{\partial z_{u}}=\frac{\partial}{\partial w_{u}}+\sum_{v,|\alpha|=\nu} C_{v, \alpha} \alpha(u) z^{\alpha_{u}} \frac{\partial}{\partial w_{v}}
$$

where $\alpha_{u}$ is the multi-index defined by the identity $z_{u} z^{\alpha_{u}}=z^{\alpha}$ (and which is undefined if $\alpha(u)=0$.) This implies immediately that:

$$
g_{u \bar{v}}(Z, G)=g_{u \bar{v}}(W, G)+\sum_{|\alpha|=n} C_{v, \alpha} \alpha(u) z^{\alpha_{u}}+\text { terms in } \bar{z}+O\left(z^{\nu+1}\right)
$$

The $\partial / \partial z_{u}$ and $\partial / \partial w_{u}$ agree to first order. Consequently:

$$
g_{u \bar{v} / \alpha_{u}}(Z, G)=g_{u \bar{v} / \alpha_{u}}(W, G)+\alpha!C_{v, \alpha}+O(z, \bar{z})
$$

Therefore the symmetric derivatives are given by:

$$
g(\alpha ; \bar{v})(Z, G)\left(z_{0}\right)=g(\alpha ; \bar{v})(W, G)\left(z_{0}\right)+\alpha!C_{v, \alpha}
$$

The identity $g(\alpha ; \bar{v})(Z, G)\left(z_{0}\right)=0$ determines the $C_{v, \alpha}$ uniquely. We take complex conjugate to conclude $g(v, \bar{\alpha})(Z, G)\left(z_{0}\right)=0$ as well.

This permits us to choose the coordinate system so all the purely holomorphic and anti-holomorphic derivatives vanish at a single point. We restrict henceforth to variables $g(\alpha ; \bar{\beta})$ so that $|\alpha| \geq 2,|\beta| \geq 2$.

We introduced the notation $h_{p \bar{q} / \alpha \bar{\beta}}$ for the jets of the metric $H$ on the auxilary coefficient bundle $V$.
Lemma 3.7.2. Let $H$ be a fiber metric on a holomorphic bundle $V$. Then given $v \geq 1$ there exists a holomorphic frame $s$ near $z_{0}$ for $V$ so that:
(a) $h_{p \bar{q}}(s, H)\left(z_{0}\right)=\delta_{p q}$.
(b) $h_{p \bar{q} / \alpha}(s, H)\left(z_{0}\right)=h_{p \bar{q} / \bar{\alpha}}(s, H)\left(z_{0}\right)=0$ for $1 \leq|\alpha| \leq \nu$.
(c) The choice of $s$ is unique modulo the action of the unitary group $\mathrm{U}(k)$
(where $k$ is the fiber dimension of $V$ ) and modulo transformations of order $\nu+1$.

Proof: If $\nu=0$, we make a linear change to assume (a). We proceed by induction assuming $s^{\prime}$ chosen for $\nu-1$. We define:

$$
s_{p}=s_{p}^{\prime}+\sum_{q,|\alpha|=\nu} C_{q, \alpha} s_{q}^{\prime} z^{\alpha}
$$

We adopt the notational conventions that indices $(p, q)$ run from 1 through $k$ and index a frame for $V$. We compute:

$$
\begin{aligned}
h_{p \bar{q}}(s, H)= & h_{p \bar{q}}\left(s^{\prime}, H\right)+\sum_{|\alpha|=\nu} C_{q, \alpha} z^{\alpha} \\
& \quad+\text { terms in } \bar{z}+\text { terms vanishing to order } \nu+1 \\
h_{p \bar{q} / \alpha}(x, H)\left(z_{0}\right)= & h_{p \bar{q} / \alpha}\left(s^{\prime}, H\right)\left(z_{0}\right)+\alpha!C_{q, \alpha} \quad \text { for }|\alpha|=\nu
\end{aligned}
$$

and derivatives of lower order are not disturbed. This determines the $C_{q, \alpha}$ uniquely so $h_{p \bar{q} / \alpha}(s, H)\left(z_{0}\right)=0$ and taking complex conjugate yields $h_{q \bar{p} / \bar{\alpha}}(s, H)\left(z_{0}\right)=0$ which completes the proof of the lemma.

Using these two normalizing lemmas, we restrict henceforth to polynomials in the variables $\left\{g(\alpha ; \bar{\beta}), h_{p \bar{q} / \alpha_{1} \bar{\beta}_{1}}\right\}$ where $|\alpha| \geq 2,|\beta| \geq 2,\left|\alpha_{1}\right| \geq 1$, $\left|\beta_{1}\right| \geq 1$. We also restrict to unitary transformations of the coordinate system and of the fiber of $V$.

If $P$ is a polynomial in these variables and if $A$ is a monomial, let $c(A, P)$ be the coefficient of $A$ in $P$. $A$ is a monomial of $P$ if $c(A, P) \neq 0$. Lemma 2.5.1 exploited invariance under the group $\mathrm{SO}(2)$ and was central to our axiomatic characterization of real Pontrjagin forms. The natural groups to study here are $\mathrm{U}(1), \mathrm{SU}(2)$ and the coordinate permutations. If we set $\partial / \partial w_{1}=a \partial / \partial z_{1}, \partial / \partial \bar{w}_{1}=\bar{a} \partial / \partial \bar{z}_{1}$ for $a \bar{a}=1$ and if we leave the other indices unchanged, then we compute:

$$
A(W, *)=a^{\operatorname{deg}_{1}(A)} \bar{a}^{\operatorname{deg}_{\overline{1}}(A)} A(Z, *) .
$$

If $A$ is a monomial of an invariant polynomial $P$, then $\operatorname{deg}_{1}(A)=\operatorname{deg}_{\overline{1}}(A)$ follows from this identity and consequently $\operatorname{deg}_{u}(A)=\operatorname{deg}_{\bar{u}}(A)$ for all $1 \leq u \leq n$.

This is the only conclusion which follows from $\mathrm{U}(1)$ invariance so we now study the group $\mathrm{SU}(2)$. We consider the coordinate transformation:

$$
\begin{array}{cc}
\partial / \partial w_{1}=a \partial / \partial z_{1}+b \partial / \partial z_{2}, & \partial / \partial w_{2}=-\bar{b} \partial / \partial z_{1}+\bar{a} \partial / \partial z_{2}, \\
\partial / \partial \bar{w}_{1}=\bar{a} \partial / \partial \bar{z}_{1}+\bar{b} \partial / \partial \bar{z}_{2}, \quad & \partial / \partial \bar{w}_{2}=-b \partial / \partial \bar{z}_{1}+a \partial / \partial \bar{z}_{2}, \\
\partial / \partial w_{u}=\partial / \partial z_{u} & \text { for } u>2, \\
\partial / \partial \bar{w}_{u}=\partial / \partial \bar{z}_{u} \quad \text { for } u>2, \\
a \bar{a}+b \bar{b}=1 .
\end{array}
$$

We let $j=\operatorname{deg}_{1}(A)+\operatorname{deg}_{2}(A)=\operatorname{deg}_{\overline{1}}(A)+\operatorname{deg}_{\overline{2}}(A)$ and expand

$$
\begin{gathered}
A(W, *)=a^{j} \bar{a}^{j} A(Z, *)+a^{j-1} \bar{a}^{j} b A(1 \rightarrow 2 \text { or } \overline{2} \rightarrow \overline{1})(Z, *) \\
+a^{j} \bar{a}^{j-1} \bar{b} A(2 \rightarrow 1 \text { or } \overline{1} \rightarrow \overline{2})(Z, *) \\
+ \text { other terms. }
\end{gathered}
$$

The notation " $A(1 \rightarrow 2$ or $\overline{2} \rightarrow \overline{1})$ " indicates all the monomials (with multiplicity) of this polynomial constructed by either changing a single index $1 \rightarrow 2$ or a single index $\overline{2} \rightarrow \overline{1}$. This plays the same role as the polynomial $A^{(1)}$ of section 2.5 where we must now also consider holomorphic and anti-holomorphic indices.

Let $P$ be $\mathrm{U}(2)$ invariant. Without loss of generality, we may assume $\operatorname{deg}_{1}(A)+\operatorname{deg}_{2}(A)=j$ is constant for all monomials $A$ of $P$ since this condition is $\mathrm{U}(2)$ invariant. (Of course, the use of the indices 1 and 2 is for notational convenience only as similar statements will hold true for any pairs of indices). We expand:

$$
\begin{aligned}
& P(W, *)=a^{j} \bar{a}^{j} P(Z, *)+a^{j-1} \bar{a}^{j} b P(1 \rightarrow 2 \text { or } \overline{2} \rightarrow \overline{1}) \\
& +a^{j} \bar{a}^{j-1} \bar{b} P(2 \rightarrow 1 \text { or } \overline{1} \rightarrow \overline{2})+O(b, \bar{b})^{2} .
\end{aligned}
$$

Since $P$ is invariant, we conclude:

$$
P(1 \rightarrow 2 \text { or } \overline{2} \rightarrow \overline{1})=P(2 \rightarrow 1 \text { or } \overline{1} \rightarrow \overline{2})=0 .
$$

We can now study these relations. Let $B$ be an arbitrary monomial and expand:

$$
B(1 \rightarrow 2)=c_{0} A_{0}+\cdots+c_{k} A_{k} \quad B(\overline{2} \rightarrow \overline{1})=-\left(d_{0} A_{0}^{\prime}+\cdots+d_{k} A_{k}^{\prime}\right)
$$

where the $c$ 's and $d$ 's are positive integers with $\sum c_{\nu}=\operatorname{deg}_{1}(B)$ and $\sum d_{\nu}=$ $\operatorname{deg}_{\overline{2}}(B)$. Then it is immediate that the $\left\{A_{j}\right\}$ and the $\left\{A_{j}^{\prime}\right\}$ denote disjoint collections of monomials since $\operatorname{deg}_{1}(A)_{\nu}=\operatorname{deg}_{1}(B)-1$ while $\operatorname{deg}_{1}(A)_{\nu}^{\prime}=$ $\operatorname{deg}_{1}(B)$. We compute:

$$
A_{\nu}(2 \rightarrow 1)=-c_{\nu}^{\prime} B+\text { other terms } \quad A_{\nu}^{\prime}(\overline{1} \rightarrow \overline{2})=d_{\nu}^{\prime} B+\text { other terms }
$$

where again the $c$ 's and $d$ 's are positive integers (related to certain multiplicities). Since $P(2 \rightarrow 1$ or $\overline{1} \rightarrow \overline{2})$ is zero, if $P$ is invariant, we conclude an identity:

$$
c(B, P(2 \rightarrow 1 \text { or } \overline{1} \rightarrow \overline{2}))=\sum_{\nu}-c_{\nu}^{\prime} c\left(A_{\nu}, P\right)+\sum_{\nu} d_{\nu}^{\prime} c\left(A_{\nu}^{\prime}, P\right)=0 .
$$

By varying the creating monomial $B$ (and also interchanging the indices 1 and 2), we can construct many linear equations among the coefficients. We note that in practice $B$ will never be a monomial of $P$ since $\operatorname{deg}_{1}(B) \neq$ $\operatorname{deg}_{\overline{1}}(B)$. We use this principle to prove the following generalization of Lemmas 2.5.1 and 2.5.5:

Lemma 3.7.3. Let $P$ be invariant under the action of $\mathrm{U}(2)$ and let $A$ be a monomial of $P$.
(a) If $A=g(\alpha ; \bar{\beta}) A_{0}^{\prime}$, then by changing only 1 and 2 indices and $\overline{1}$ and $\overline{2}$ indices we can construct a new monomial $A_{1}$ of $P$ which has the form $A_{1}=g\left(\alpha_{1} ; \bar{\beta}_{1}\right) A_{0}^{\prime \prime}$ where $\alpha(1)+\alpha(2)=\alpha_{1}(1)$ and where $\alpha_{1}(2)=0$.
(b) If $A=h_{p \bar{q} / \alpha \bar{\beta}} A_{0}^{\prime}$ then by changing only 1 and 2 indices and $\overline{1}$ and $\overline{2}$ indices we can construct a new monomial $A_{1}$ of $P$ which has the form $A_{1}=h_{p \bar{q} / \alpha_{1} \bar{\beta}_{1}} A_{0}^{\prime \prime}$ where $\alpha(1)+\alpha(2)=\alpha_{1}(1)$ and where $\alpha_{1}(2)=0$.
Proof: We prove (a) as the proof of (b) is the same. Choose $A$ of this form so $\alpha(1)$ is maximal. If $\alpha(2)$ is zero, we are done. Suppose the contrary. Let $B=g\left(\alpha_{1} ; \bar{\beta}\right) A_{0}^{\prime}$ where $\alpha_{1}=(\alpha(1)+1, \alpha(2)-1, \alpha(3), \ldots, \alpha(n))$. We expand:

$$
\begin{aligned}
& B(1 \rightarrow 2)=\alpha_{1}(1) A+\text { monomials divisible by } g\left(\alpha_{1} ; \bar{\beta}\right) \\
& B(\overline{2} \rightarrow \overline{1})=\text { monomials divisible by } g\left(\alpha_{1} ; \bar{\beta}_{1}\right) \text { for some } \beta_{1} .
\end{aligned}
$$

Since $A$ is a monomial of $P$, we use the principle described previously to conclude there is some monomial of $P$ divisible by $g\left(\alpha_{1} ; \bar{\beta}_{1}\right)$ for some $\beta_{1}$. This contradicts the maximality of $\alpha$ and shows $\alpha(2)=0$ completing the proof.

We now begin the proof of Theorem 3.6.9. Let $0 \neq P \in \mathcal{R}_{n, \nu, k}^{c}$ be a scalar valued invariant homogeneous of order $\nu$ with $r(P)=0$. Let $A$ be a monomial of $P$. Decompose

$$
A=g\left(\alpha_{1} ; \bar{\beta}_{1}\right) \ldots g\left(\alpha_{r} ; \bar{\beta}_{r}\right) h_{p_{1} q_{1} / \alpha_{r+1} \bar{\beta}_{r+1}} \ldots h_{p_{s} \bar{q}_{s} / \alpha_{r+s} \bar{\beta}_{r+s}} .
$$

Let $\ell(A)=r+s$ be the length of $A$. We show $\ell(A) \geq n$ as follows. Without changing $\left(r, s,\left|\alpha_{\nu}\right|,\left|\beta_{\nu}\right|\right)$ we can choose $A$ in the same form so that $\alpha_{1}(u)=0$ for $u>1$. We now fix the index 1 and apply Lemma 3.7.3 to the remaining indices to choose $A$ so $\alpha_{2}(u)=0$ for $\nu>2$. We continue in this fashion to construct such an $A$ so that $\operatorname{deg}_{u}(A)=0$ for $u>r+s$. Since $\operatorname{deg}_{n}(A)=\operatorname{deg}_{\bar{n}}(A)>0$, we conclude therefore that $r+s \geq n$. We estimate:

$$
\begin{aligned}
\nu=\operatorname{ord}(P) & =\sum_{\mu \leq r}\left\{\left|\alpha_{\mu}\right|+\left|\beta_{\mu}\right|-2\right\}+\sum_{r<\mu \leq r+s}\left\{\left|\alpha_{\mu}\right|+\left|\beta_{u}\right|\right\} \\
& \geq 2 r+2 s \geq 2 n=m
\end{aligned}
$$

since $\operatorname{ord}(g(\alpha, \beta)) \geq 2$ and $\operatorname{ord}\left(h_{p \bar{q} / \alpha \bar{\beta}}\right) \geq 2$. This proves $\nu \geq m \geq 2$. Consequently, if $\nu<m$ we conclude $P=0$ which proves the first assertion of Theorem 3.6.

We assume henceforth that we are in the limiting case $\nu=m$. In this case, all the inequalities must have been equalities. This implies $r+s=n$ and

$$
\left|\alpha_{\mu}\right|+\left|\beta_{\mu}\right|=4 \quad \text { for } \mu \leq r \quad \text { and } \quad\left|\alpha_{\mu}\right|+\left|\beta_{\mu}\right|=2 \quad \text { for } r<\mu \leq s+r .
$$

This holds true for every monomial of $P$ since the construction of $A_{1}$ from $A$ involved the use of Lemma 3.7.3 and does not change any of the orders involved. By Lemma 3.7.1 and 3.7.2 we assumed all the purely holomorphic and purely anti-holomorphic derivatives vanished at $z_{0}$ and consequently:

$$
\left|\alpha_{\mu}\right|=\left|\beta_{\mu}\right|=2 \quad \text { for } \mu \leq r \quad \text { and } \quad\left|\alpha_{\mu}\right|=\left|\beta_{\mu}\right|=1 \quad \text { for } r<\mu \leq s+r .
$$

Consequently, $P$ is a polynomial in the $\left\{g\left(i_{1} i_{2} ; \bar{\jmath}_{1} \bar{\jmath}_{2}\right) h_{p \bar{q} / i \bar{\jmath}}\right\}$ variables; it only involves the mixed 2-jets involved.

We wish to choose a monomial of $P$ in normal form to begin counting the number of possible such $P$. We begin this process with:
Lemma 3.7.4. Let $P$ satisfy the hypothesis of Theorem 3.6.8(b). Then there exists a monomial $A$ of $P$ which has the form:

$$
A=g\left(11 ; \bar{\jmath}_{1} \jmath_{1}^{\prime}\right) \ldots g\left(t t ; \bar{\jmath}_{t} \bar{\jmath}_{t}^{\prime}\right) h_{p_{1} q_{1} / t+1, \bar{\jmath}_{t+1}} \ldots h_{p_{s} q_{s} / m, \bar{\jmath}_{n}}
$$

where $t+s=n$.
Proof: We let $*$ denote indices which are otherwise unspecified and let $A_{0}$ be a generic monomial. We have $\operatorname{deg}_{n}(A)>0$ for every monomial $A$ of $P$ so $\operatorname{deg}_{u}(A)>0$ for every index $u$ as well. We apply Lemma 3.7.3 to choose $A$ of the form:

$$
A=g\left(11 ; \bar{\jmath}_{1} \bar{\jmath}_{1}^{\prime}\right) A_{0} .
$$

Suppose $r>1$. If $A=g(11 ; *) g(11: *) A_{0}$, then we could argue as before using Lemma 3.7.3 that we could choose a monomial $A$ of $P \operatorname{sodeg}_{k}(A)=0$ for $k>1+\ell\left(A_{0}\right)=r+s-1=n-1$ which would be false. Thus $A=$ $g(11 ; *) g(j k ; *) A_{0}$ where not both $j$ and $k$ are 1 . We may apply a coordinate permutation to choose $A$ in the form $A=g(11 ; *) g(2 j ; *) A_{0}$. If $j \geq 2$ we can apply Lemma 3.7 .3 to choose $A=g(11 ; *) g(22 ; *) A_{0}$. Otherwise we suppose $A=g(11 ; *) g(12 ; *) A_{0}$. We let $B=g(11 ; *) g(22 ; *) A_{0}$ and compute:

$$
\begin{aligned}
& B(2 \rightarrow 1)=A+\text { terms divisible by } g(11 ; *) g(22 ; *) \\
& B(\overline{1} \rightarrow \overline{2})=\text { terms divisible by } g(11 ; *) g(22 ; *)
\end{aligned}
$$

so that we conclude in any event we can choose $A=g(11 ; *) g(22 ; *) A_{0}$. We continue this argument with the remaining indices to construct:

$$
A=g(11 ; *) g(22 ; *) \ldots g(t t ; *) h_{p_{1} \bar{q}_{1} / u_{1} \bar{v}_{1}} \ldots h_{p_{s} \bar{q}_{s} / u_{s} \bar{v}_{s}} .
$$

We have $\operatorname{deg}_{u}(A) \neq 0$ for all $u$. Consequently, the indices $\{t+1, \ldots, t+s\}$ must appear among the indices $\left\{u_{1}, \ldots, u_{s}\right\}$. Since these two sets have $s$ elements, they must coincide. Thus by rearranging the indices we can assume $u_{\nu}=\nu+t$ which completes the proof. We note $\operatorname{deg}_{u}(A)=2$ for $u \leq t$ and $\operatorname{deg}_{u}(A)=1$ for $u>t$.

This lemma does not control the anti-holomorphic indices, we further normalize the choice of $A$ in the following:

Lemma 3.7.5. Let $P$ satisfy the hypothesis of Theorem 3.6.8(b). Then there exists a monomial $A$ of $P$ which has the form:

$$
A=g\left(11 ; \bar{\jmath}_{1} \bar{\jmath}_{1}^{\prime}\right) \ldots g\left(t t ; \bar{\jmath}_{t} \bar{\jmath}_{t}^{\prime}\right) h_{p_{1} \bar{q}_{1} / t+1, \overline{t+1}} \ldots h_{p_{s} \bar{q}_{s} / n \bar{n}}
$$

for

$$
1 \leq j_{\mu}, j_{\mu}^{\prime} \leq t
$$

Proof: If $A$ has this form, then $\operatorname{deg}_{j}(A)=1$ for $j>t$. Therefore $\operatorname{deg}_{\bar{\jmath}}(A)=1$ for $j>t$ which implies $1 \leq j_{\mu}, j_{\mu}^{\prime} \leq t$ automatically. We say that $\bar{\jmath}$ touches itself in $A$ if $A$ is divisible by $g(* ; \bar{\jmath} \bar{\jmath})$ for some $*$. We say that $j$ touches $\bar{\jmath}$ in $A$ if $A$ is divisible by $h_{p q / j \bar{\jmath}}$ for some $(p, q)$. Choose $A$ of the form given in Lemma 3.7.4 so the number of indices $j>t$ which touch $\bar{\jmath}$ in $A$ is maximal. Among all such $A$, choose $A$ so the number of $\bar{\jmath} \leq \bar{t}$ which touch themselves in $A$ is maximal. Suppose $A$ does not satisfy the conditions of the Lemma. Thus $A$ must be divisible by $h_{p \bar{q} / u \bar{v}}$ for $u \neq v$ and $u>t$. We suppose first $v>t$. Since $\operatorname{deg}_{\bar{v}}(A)=\operatorname{deg}_{v}(A)=1$, $v$ does not touch $\bar{v}$ in $A$. We let $A=h_{p \bar{q} / u \bar{v}} A_{0}$ and $B=h_{p \bar{q} / u \bar{u}} A_{0}$ then $\operatorname{deg}_{\bar{v}}(B)=0$. We compute:

$$
B(\bar{u} \rightarrow \bar{v})=A+A_{1} \quad \text { and } \quad B(v \rightarrow u)=A_{2}
$$

where $A_{1}$ is defined by interchanging $\bar{u}$ and $\bar{v}$ and where $A_{2}$ is defined by replacing both $v$ and $\bar{v}$ by $u$ and $\bar{u}$ in $A$. Thus $\operatorname{deg}_{v}(A)_{2}=0$ so $A_{2}$ is not a monomial of $A$. Thus $A_{1}=h_{p \bar{q} / u \bar{u}} A_{0}^{\prime}$ must be a monomial of $P$. One more index (namely $u$ ) touches its holomorphic conjugate in $A_{1}$ than in $A$. This contradicts the maximality of $A$ and consequently $A=h_{p \bar{q} / u \bar{v}} A_{0}^{\prime}$ for $u>t, v \leq t$. (This shows $h_{p_{0} q_{0} / u_{0} \bar{v}_{0}}$ does not divide $A$ for any $u_{0} \neq v_{0}$ and $\bar{v}_{0}>\bar{t}$.) We have $\operatorname{deg}_{\bar{u}}(A)=1$ so the anti-holomorphic index $\bar{u}$ must appear somewhere in $A$. It cannot appear in a $h$ variable and consequently $A$ has the form $A=g(* ; \bar{u} \bar{w}) h_{p \bar{q} / u \bar{v}} A_{0}^{\prime \prime}$. We define $B=g(* ; \bar{w} \bar{w}) h_{p \bar{q} / u \bar{v}} A_{0}^{\prime \prime}$ and compute:

$$
\begin{aligned}
& B(\bar{w} \rightarrow \bar{u})=A+\text { terms divisible by } g(* ; \bar{w} \bar{w}) \\
& B(u \rightarrow w)=g(* ; \bar{w} \bar{w}) h_{p \bar{q} / w \bar{v}} A_{0}^{\prime \prime} .
\end{aligned}
$$

Since $g(* ; \bar{w} \bar{w}) h_{p \bar{q} / w \bar{v}} A_{0}^{\prime \prime}$ does not have the index $u$, it cannot be a monomial of $P$. Thus terms divisible by $g(* ; \bar{w} \bar{w})$ must appear in $P$. We construct these terms by interchanging a $\bar{w}$ and $\bar{u}$ index in $P$ so the maximality of indices $u_{0}$ touching $\bar{u}_{0}$ for $u_{0}>t$ is unchanged. Since $\bar{w}$ does not touch $\bar{w}$ in $A$, we are adding one additional index of this form which again contradicts the maximality of $A$. This final contradiction completes the proof.

This constructs a monomial $A$ of $P$ which has the form

$$
A=A_{0} A_{1} \quad \text { for } \quad \begin{cases}\operatorname{deg}_{u}\left(A_{0}\right)=\operatorname{deg}_{\bar{u}}\left(A_{0}\right)=0, & u>t \\ \operatorname{deg}_{u}\left(A_{1}\right)=\operatorname{deg}_{\bar{u}}\left(A_{1}\right)=0, & u \leq t\end{cases}
$$

$A_{0}$ involves only the derivatives of the metric $g$ and $A_{1}$ involves only the derivatives of $h$. We use this splitting in exactly the same way we used a similar splitting in the proof of Theorem 2.6.1 to reduce the proof of Theorem 3.6.9 to the following assertions:
Lemma 3.7.6. Let $P$ satisfy the hypothesis of Theorem 3.6.9(b).
(a) If $P$ is a polynomial in the $\left\{h_{p \bar{q} / u \bar{v}}\right\}$ variables, then $P=* Q$ for $Q$ a Chern $m$ form of $V$.
(b) If $P$ is a polynomial in the $\left\{g\left(u_{1} u_{2} ; \bar{v}_{1} \bar{v}_{2}\right)\right\}$ variables, then $P=* Q$ for $Q$ a Chern $m$ form of $T_{c} M$.
Proof: (a). If $A$ is a monomial of $P$, then $\operatorname{deg}_{u}(A)=\operatorname{deg}_{\bar{u}}(A)=1$ for all $u$ so we can express

$$
A=h_{p_{1} q_{1} / 1 \bar{u}_{1}} \ldots h_{p_{n} q_{n} / n \bar{u}_{n}}
$$

where the $\left\{u_{\nu}\right\}$ are a permutation of the indices $i$ through $n$. We let $A=h_{* / 1 \bar{u}_{1}} h_{* / 2 \bar{u}_{2}} A_{0}^{\prime}$ and $B=h_{* / 1 \bar{u}_{1}} h_{* / 2 \bar{u}_{1}} A_{0}^{\prime}$. Then:

$$
\begin{array}{ll}
B\left(\bar{u}_{1} \rightarrow \bar{u}_{2}\right)=A+A_{1} & \text { for } A_{1}=h_{* / 1 \bar{u}_{2}} h_{* / 2 \bar{u}_{1}} A_{0}^{\prime} \\
B\left(u_{2} \rightarrow u_{1}\right)=A_{2} & \text { for } \operatorname{deg}_{u_{2}}\left(A_{2}\right)=0 .
\end{array}
$$

Therefore $A_{2}$ is not a monomial of $P$. Thus $A_{1}$ is a monomial of $P$ and furthermore $c(A, P)+c\left(A_{1}, P\right)=0$. This implies when we interchange $\bar{u}_{1}$ and $\bar{u}_{2}$ that we change the sign of the coefficient involved. This implies immediately we can express $P$ is terms of expressions:

$$
\begin{aligned}
*\left(\Omega_{p_{1} \bar{q}_{1}}^{V} \wedge \cdots \wedge \Omega_{p_{n} \bar{q}_{n}}^{V}\right)= & \sum_{\vec{u}, \vec{v}} *\left(h_{p_{1} \bar{q}_{1} / u_{1} \bar{v}_{1}} d u^{1} \wedge d \bar{v}_{1} \cdots\right. \\
& \left.\wedge h_{p_{n} \bar{q}_{n} / u_{n} \bar{v}_{n}} d u^{n} \wedge d v^{n}\right) \\
= & n!\sum_{\rho} \operatorname{sign}(\rho) h_{p_{1} q_{1} / 1 \bar{\rho}(1)} \ldots h_{p_{n} q_{n} / n \bar{\rho}(n)}+\cdots
\end{aligned}
$$

where $\rho$ is a permutation. This implies $* P$ can be expressed as an invariant polynomial in terms of curvature which implies it must be a Chern form as previously computed.

The remainder of this section is devoted to the proof of (b).
Lemma 3.7.7. Let $P$ satisfy the hypothesis of Lemma 3.7.6(b). Then we can choose a monomial $A$ of $P$ which has the form:

$$
A=g\left(11 ; \bar{u}_{1} \bar{u}_{1}\right) \ldots g\left(n n ; \bar{u}_{n} \bar{u}_{n}\right) .
$$

This gives a normal form for a monomial. Before proving Lemma 3.7.7, we use this lemma to complete the proof of Lemma 3.7.6(b). By making a
coordinate permutation if necessary we can assume $A$ has either the form $g(11 ; \overline{1} \overline{1}) A_{0}^{\prime}$ or $g(11 ; \overline{2} \overline{2}) A_{0}^{\prime}$. In the latter case, we continue inductively to express $A=g(11 ; \overline{2} \overline{2}) g(22 ; \overline{3} \overline{3}) \ldots g(u-1, u-1 ; \bar{u} \bar{u}) g(u u ; \overline{1} \overline{1}) A_{0}^{\prime}$ until the cycle closes. If we permit $u=1$ in this decomposition, we can also include the first case. Since the indices 1 through $u$ appear exactly twice in $A$ they do not appear in $A_{0}^{\prime}$. Thus we can continue to play the same game to decompose $A$ into cycles. Clearly $A$ is determined by the length of the cycles involved (up to coordinate permutations); the number of such classifying monomials is $\pi(n)$, the number of partitions of $n$. This shows that the dimension of the space of polynomials $P$ satisfying the hypothesis of Lemma $3.7 .6(\mathrm{~b})$ is $\leq \pi(n)$. Since there are exactly $\pi(n)$ Chern forms the dimension must be exactly $\pi(n)$ and every such $P$ must be a Chern $m$ form as claimed.

We give an indirect proof to complete the proof of Lemma 3.7.7. Choose $A$ of the form given by Lemma 3.7.5 so the number of anti-holomorphic indices which touch themselves is maximal. If every anti-holomorphic index touches itself, then $A$ has the form of Lemma 3.7.7 and we are done. We suppose the contrary. Since every index appears exactly twice, every antiholomorphic index which does not touch itself touches another index which also does not touch itself. Every holomorphic index touches only itself. We may choose the notation so $A=g(* ; \overline{1} \overline{2}) A_{0}^{\prime}$. Suppose first $\overline{1}$ does not touch $\overline{2}$ in $A_{0}^{\prime}$. Then we can assume $A$ has the form:

$$
A=g(* ; \overline{1} \overline{2}) g(* ; \overline{1} \overline{3}) g(* ; \overline{2} \bar{k}) A_{0}^{\prime}
$$

where possibly $k=3$ in this expression. The index 1 touches itself in $A$. The generic case will be:

$$
A=g(11 ; *) g(* ; \overline{1} \overline{2}) g(* ; \overline{1} \overline{3}) g(* ; \overline{2} \bar{k}) A_{0}^{\prime} .
$$

The other cases in which perhaps $A=g(11 ; \overline{1} \overline{2}) \ldots$ or $g(11, \overline{1} \overline{3}) \ldots$ or $g(11 ; \overline{2} \bar{k}) \ldots$ are handled similarly. The holomorphic and anti-holomorphic indices do not interact. In exactly which variable they appear does not matter. This can also be expressed as a lemma in tensor algebras.

We suppose $k \neq 3$; we will never let the 2 and 3 variables interact so the case in which $k=3$ is exactly analogous. Thus $A$ has the form:

$$
A=g(11 ; *) g(* ; \overline{1} \overline{2}) g(* ; \overline{1} \overline{3}) g(* ; \overline{2} \bar{k}) g(* ; \overline{3} \bar{\jmath}) A_{0}^{\prime}
$$

where possibly $j=k$. Set $B=(11 ; *) g(* ; \overline{2} \overline{2}) g(* ; \overline{1} \overline{3}) g(* ; \overline{2} \bar{k}) g(* ; \overline{3} \bar{\jmath}) A_{0}^{\prime}$ then:

$$
\begin{aligned}
& B(\overline{2} \rightarrow \overline{1})=2 A+A^{\prime \prime} \\
& B(1 \rightarrow 2)=2 A_{1}
\end{aligned}
$$

for

$$
\begin{aligned}
& A^{\prime \prime}=g(11 ; *) g(* ; \overline{2} \overline{2}) g(* ; \overline{1} \overline{3}) g(* ; \overline{1} \bar{k}) g(* ; \overline{3} \bar{\jmath}) A_{0}^{\prime} \\
& A_{1}=g(12 ; *) g(* ; \overline{2} \overline{2}) g(* ; \overline{1} \overline{3}) g(* ; \overline{2} \bar{k}) g(* ; \overline{3} \bar{\jmath}) A_{0}^{\prime}
\end{aligned}
$$

$A^{\prime \prime}$ is also a monomial of the form given by Lemma 3.7.5. Since one more anti-holomorphic index touches itself in $A^{\prime \prime}$ then does in $A$, the maximality of $A$ shows $A^{\prime \prime}$ is not a monomial of $P$. Consequently $A_{1}$ is a monomial of $P$. Set $B_{1}=g(12 ; *) g(* ; \overline{2} \overline{2}) g(* ; \overline{3} \overline{3}) g(* ; \overline{2} \bar{k}) g(* ; \overline{3} \bar{\jmath}) A_{0}^{\prime}$ then:

$$
\begin{aligned}
& B_{1}(\overline{3} \rightarrow \overline{1})=2 A_{1}+A_{2} \\
& B_{1}(1 \rightarrow 3)=A^{\prime \prime \prime}
\end{aligned}
$$

for

$$
\begin{aligned}
A_{2} & =g(12 ; *) g(* ; \overline{2} \overline{2}) g(* ; \overline{3} \overline{3}) g(* ; \overline{2} \bar{k}) g(* ; \overline{1} \bar{\jmath}) A_{0}^{\prime} \\
A^{\prime \prime \prime} & =g(32 ; *) g(* ; \overline{2} \overline{2}) g(* ; \overline{3} \overline{3}) g(* ; \overline{2} \bar{k}) g(* ; \overline{3} \bar{\jmath}) A_{0}^{\prime}
\end{aligned}
$$

However, $\operatorname{deg}_{1}\left(A^{\prime \prime \prime}\right)=0$. Since $r(P)=0, A^{\prime \prime \prime}$ cannot be a monomial of $P$ so $A_{2}$ is a monomial of $P$. Finally we set

$$
B_{2}=g(11 ; *) g(* ; \overline{2} \overline{2}) g(* ; \overline{3} \overline{3}) g(* ; \overline{2} \bar{k}) g(* ; \overline{1} \bar{\jmath}) A_{0}^{\prime}
$$

so

$$
\begin{aligned}
& B_{2}(1 \rightarrow 2)=A_{2} \\
& B_{2}(\overline{2} \rightarrow \overline{1})=A_{3}+2 A_{4}
\end{aligned}
$$

for

$$
\begin{aligned}
& A_{3}=g(11 ; *) g(* ; \overline{2} \overline{2}) g(* ; \overline{3} \overline{3}) g(* ; \overline{1} \bar{k}) g(* ; \overline{1} \bar{\jmath}) A_{0}^{\prime} \\
& A_{4}=g(11 ; *) g(* ; \overline{1} \overline{2}) g(* ; \overline{3} \overline{3}) g(* ; \overline{2} \bar{k}) g(* ; \overline{1} \bar{\jmath}) A_{0}^{\prime}
\end{aligned}
$$

This implies either $A_{3}$ or $A_{4}$ is monomial of $P$. Both these have every holomorphic index touching itself. Furthermore, one more anti-holomorphic index (namely $\overline{3}$ ) touches itself. This contradicts the maximality of $A$.

In this argument it was very important that $\overline{2} \neq \overline{3}$ as we let the index $\overline{1}$ interact with each of these indices separately. Thus the final case we must consider is the case in which $A$ has the form:

$$
A=g(11 ; *) g(22 ; *) g(* ; \overline{1} \overline{2}) g(* ; \overline{1} \overline{2}) A_{0} .
$$

So far we have not had to take into account multiplicities or signs in computing $A(1 \rightarrow 2)$ etc; we have been content to conclude certain coefficients are non-zero. In studying this case, we must be more careful in our analysis as the signs involved are crucial. We clear the previous notation and define:

$$
\begin{aligned}
& A_{1}=g(12 ; *) g(22 ; *) g(* ; \overline{1} \overline{2}) g(* ; \overline{2} \overline{2}) A_{0} \\
& A_{2}=g(12 ; *) g(22 ; *) g(* ; \overline{2} \overline{2}) g(* ; \overline{1} \overline{2}) A_{0} \\
& A_{3}=g(11 ; *) g(22 ; *) g(* ; \overline{1} \overline{1}) g(* ; \overline{2} \overline{2}) A_{0} \\
& A_{4}=g(22 ; *) g(22 ; *) g(* ; \overline{2} \overline{2}) g(* ; \overline{2} \overline{2}) A_{0}
\end{aligned}
$$

We note that $A_{3}$ is not a monomial of $P$ by the maximality of $A . A_{4}$ is not a monomial of $P$ as $\operatorname{deg}_{1}\left(A_{4}\right)=0$ and $r(P)=0$. We let $B=$ $g(11 ; *) g(22 ; *) g(* ; \overline{1} \overline{2}) g(* ; \overline{2} \overline{2}) A_{0}$. Then

$$
B(\overline{2} \rightarrow \overline{1})=2 A+A_{3}, \quad B(1 \rightarrow 2)=2 A_{1}
$$

so that $A_{1}$ must be a monomial of $P$ since $A$ is a monomial of $P$ and $A_{4}$ is not. We now pay more careful attention to the multiplicities and signs:

$$
A(\overline{1} \rightarrow \overline{2})=B+\cdots \quad \text { and } \quad A_{1}(2 \rightarrow 1)=B+\cdots
$$

However $\overline{1} \rightarrow \overline{2}$ introduces a $\bar{b}$ while $2 \rightarrow 1$ introduces a $-\bar{b}$ so using the argument discussed earlier we conclude not only $c\left(A_{1}, P\right) \neq 0$ but that $c(A, P)-c\left(A_{1}, P\right)=0$ so $c\left(A_{1}, P\right)=c(A, P) . \quad A_{2}$ behaves similarly so the analogous argument using $B_{1}=g(11 ; *) g(22 ; *) g(* ; \overline{2} \overline{2}) g(* ; \overline{1} \overline{2})$ shows $c\left(A_{2}, P\right)=c(A, P)$. We now study $B_{2}=g(12 ; *) g(22 ; *) g(*, \overline{2} \overline{2}) g(* ; \overline{2} \overline{2})$ and compute:

$$
B_{2}(1 \rightarrow 2)=A_{4} \quad B_{2}(\overline{2} \rightarrow \overline{1})=2 A_{1}+2 A_{2}
$$

$A_{4}$ is not a monomial of $P$ as noted above. We again pay careful attention to the signs:

$$
A_{1}(\overline{1} \rightarrow \overline{2})=B_{2} \quad \text { and } \quad A_{2}(\overline{1} \rightarrow \overline{2})=B_{2}
$$

This implies $c\left(A_{1}, P\right)+c\left(A_{2}, P\right)=0$. Since $c\left(A_{1}, P\right)=c\left(A_{2}, P\right)=c(A, P)$ this implies $2 c(A, P)=0$ so $A$ was not a monomial of $P$. This final contradiction completes the proof.

This proof was long and technical. However, it is not a theorem based on unitary invariance alone as the restriction axiom plays an important role in the development. We know of no proof of Theorem 3.6.9 which is based only on H. Weyl's theorem; in the real case, the corresponding characterization of the Pontrjagin classes was based only on orthogonal invariance and we gave a proof based on H. Weyl's theorem in that case.

### 3.8. The Chern Isomorphism and Bott Periodicity.

In section 3.9 we shall discuss the Atiyah-Singer index theorem in general using the results of section 3.1. The index theorem gives a topological formula for the index of an arbitrary elliptic operator. Before begining the proof of that theorem, we must first review briefly the Bott periodicity theorem from the point of Clifford modules. We continue our consideration of the bundles over $S^{n}$ constructed in the second chapter. Let

$$
\begin{aligned}
\mathrm{GL}(k, \mathbf{C}) & =\{A: A \text { is a } k \times k \text { complex matrix with } \operatorname{det}(A) \neq 0\} \\
\mathrm{GL}^{\prime}(k, \mathbf{C}) & =\{A \in \mathrm{GL}(k, \mathbf{C}): \operatorname{det}(A-i t) \neq 0 \text { for all } t \in \mathbf{R}\} \\
\mathrm{U}(k) & =\left\{A \in \mathrm{GL}(k, \mathbf{C}): A \cdot A^{*}=I\right\} \\
\mathrm{S}(k) & =\left\{A \in \mathrm{GL}(k, \mathbf{C}): A^{2}=I \text { and } A=A^{*}\right\} \\
\mathrm{S}_{0}(k) & =\{A \in \mathrm{~S}(k): \operatorname{Tr}(A)=0\} .
\end{aligned}
$$

We note that $\mathrm{S}_{0}(k)$ is empty if $k$ is odd. $\mathrm{U}(k)$ is compact and is a deformation retract of $\mathrm{GL}(k, \mathbf{C}) ; \mathrm{S}(k)$ is compact and is a deformation retract of $\mathrm{GL}^{\prime}(k, \mathbf{C}) . \mathrm{S}_{0}(k)$ is one of the components of $\mathrm{S}(k)$.

Let $X$ be a finite simplicial complex. The suspension $\Sigma X$ is defined by identifying $X \times\left\{\frac{\pi}{2}\right\}$ to a single point $N$ and $X \times\left\{-\frac{\pi}{2}\right\}$ to single point $S$ in the product $X \times\left[-\frac{\pi}{2}, \frac{\pi}{2}\right]$. Let $D_{ \pm}(X)$ denote the northern and southern "hemispheres" in the suspension; the intersection $D_{+}(X) \cap D_{-}(X)=X$.
[Reprinter's note: A figure depicting the suspension of $X$ belongs here.]

We note both $D_{+}(X)$ and $D_{-}(X)$ are contractible. $\Sigma(X)$ is a finite simplicial complex. If $S^{n}$ is the unit sphere in $\mathrm{R}^{n+1}$, then $\Sigma\left(S^{n}\right)=S^{n+1}$. Finally, if $W$ is a vector bundle over some base space $Y$, then we choose a fiber metric on $W$ and let $S(W)$ be the unit sphere bundle. $\Sigma(W)$ is the fiberwise suspension of $S(W)$ over $Y$. This can be identified with $S(W \oplus 1)$.

It is beyond the scope of this book to develop in detail the theory of vector bundles so we shall simply state relevant facts as needed. We let $\operatorname{Vect}_{k}(X)$ denote the set of isomorphism classes of complex vector bundles
over $X$ of fiber dimension $k$. We let $\operatorname{Vect}(X)=\bigcup_{k} \operatorname{Vect}_{k}(X)$ be the set of isomorphism classes of complex vector bundles over $X$ of all fiber dimensions. We assume $X$ is connected so that the dimension of a vector bundle over $X$ is constant.

There is a natural inclusion map $\operatorname{Vect}_{k}(X) \rightarrow \operatorname{Vect}_{k+1}(X)$ defined by sending $V \mapsto V \oplus 1$ where 1 denotes the trivial line bundle over $X$.

Lemma 3.8.1. If $2 k \geq \operatorname{dim} X$ then the $\operatorname{map}_{\operatorname{Vect}}^{k}(X) \rightarrow \operatorname{Vect}_{k+1}(X)$ is bijective.

We let $F(\operatorname{Vect}(X))$ be the free abelian group on these generators. We shall let $(V)$ denote the element of $K(X)$ defined by $V \in \operatorname{Vect}(X) . K(X)$ is the quotient modulo the relation $(V \oplus W)=(V)+(W)$ for $V, W \in \operatorname{Vect}(X)$. $K(X)$ is an abelian group. The natural map $\operatorname{dim}: \operatorname{Vect}(X) \rightarrow \mathbf{Z}^{+}$extends to $\operatorname{dim}: K(X) \rightarrow \mathbf{Z} ; \widetilde{K}(X)$ is the kernel of this map. We shall say that an element of $K(X)$ is a virtual bundle; $\widetilde{K}(X)$ is the subgroup of virtual bundles of virtual dimension zero.

It is possible to define a kind of inverse in $\operatorname{Vect}(X)$ :
Lemma 3.8.2. Given $V \in \operatorname{Vect}(X)$, there exists $W \in \operatorname{Vect}(X)$ so $V \oplus W=$ $1^{j}$ is isomorphic to a trivial bundle of dimension $j=\operatorname{dim} V+\operatorname{dim} W$.

We combine these two lemmas to give a group structure to $\operatorname{Vect}_{k}(X)$ for $2 k \geq \operatorname{dim} X$. Given $V, W \in \operatorname{Vect}_{k}(X)$, then $V \oplus W \in \operatorname{Vect}_{2 k}(X)$. By Lemma 3.8.1, the map $\operatorname{Vect}_{k}(X) \rightarrow \operatorname{Vect}_{2 k}(X)$ defined by sending $U \mapsto U \oplus 1^{k}$ is bijective. Thus there is a unique element we shall denote by $V+W \in$ $\operatorname{Vect}_{k}(X)$ so that $(V+W) \oplus 1^{k}=V \oplus W$. It is immediate that this is an associative and commutative operation and that the trivial bundle $1^{k}$ functions as the unit. We use Lemma 3.8.2 to construct an inverse. Given $V \in \operatorname{Vect}_{k}(X)$ there exists $j$ and $W \in \operatorname{Vect}_{j}(X)$ so $V \oplus W=1^{j+k}$. We assume without loss of generality that $j \geq k$. By Lemma 3.8.1, we choose $\bar{W} \in \operatorname{Vect}_{k}(X)$ so that $V \oplus \bar{W} \oplus 1^{j-k}=V \oplus \bar{W}=1^{j+k}$. Then the bijectivity implies $V \oplus \bar{W}=1^{2 k}$ so $\bar{W}$ is the inverse of $V$. This $\operatorname{shows}^{\operatorname{Vect}}{ }_{k}(X)$ is a group under this operation.

There is a natural map $\operatorname{Vect}_{k}(X) \rightarrow \widetilde{K}(X)$ defined by sending $V \mapsto$ $(V)-\left(1^{k}\right)$. It is immediate that:

$$
\begin{aligned}
(V+W)-1^{k} & =(V+W)+\left(1^{k}\right)-\left(1^{2 k}\right)=\left((V+W) \oplus 1^{k}\right)-\left(1^{2 k}\right) \\
& =(V \oplus W)-\left(1^{2 k}\right) \\
& =(V)+(W)-\left(1^{2 k}\right)=(V)-\left(1^{k}\right)+(W)-\left(1^{k}\right)
\end{aligned}
$$

so the map is a group homomorphism. $\widetilde{K}(X)$ is generated by elements of the form $(V)-(W)$ for $V, W \in \operatorname{Vect}_{j} X$ for some $j$. If we choose $\bar{W}$ so $W \oplus \bar{W}=1^{v}$ then $(V)-(W)=(V \oplus \bar{W})-\left(1^{v}\right)$ so $\widetilde{K}(X)$ is generated by elements of the form $(V)-\left(1^{v}\right)$ for $V \in \operatorname{Vect}_{v}(X)$. Again, by adding
trivial factors, we may assume $v \geq k$ so by Lemma 3.8.1 $V=V_{1} \oplus 1^{j-k}$ and $(V)-\left(1^{j}\right)=\left(V_{1}\right)-\left(1^{k}\right)$ for $V_{1} \in \operatorname{Vect}_{k}(X)$. This implies the map $\operatorname{Vect}_{k}(X) \rightarrow \widetilde{K}(X)$ is subjective. Finally, we note that that in fact $K(X)$ is generated by $\operatorname{Vect}_{k}(X)$ subject to the relation $(V)+(W)=(V \oplus W)=$ $(V+W)+\left(1^{k}\right)$ so that this map is injective and we may identify $\widetilde{K}(X)=$ $\operatorname{Vect}_{k}(X)$ and $K(X)=\mathbf{Z} \oplus \widetilde{K}(X)=\mathbf{Z} \oplus \operatorname{Vect}_{k}(X)$ for any $k$ such that $2 k \geq \operatorname{dim} X$.

Tensor product defines a ring structure on $K(X)$. We define $(V) \otimes(W)=$ $(V \otimes W)$ for $V, W \in \operatorname{Vect}(X)$. Since $\left(V_{1} \oplus V_{2}\right) \otimes(W)=\left(V_{1} \otimes W \oplus V_{2} \otimes W\right)$, this extends from $F(\operatorname{Vect}(X))$ to define a ring structure on $K(X)$ in which the trivial line bundle 1 functions as a multiplicative identity. $\widetilde{K}(X)$ is an ideal of $K(X)$.
$K(X)$ is a $\mathbf{Z}$-module. It is convenient to change the coefficient group and define:

$$
K(X ; \mathbf{C})=K(X) \otimes_{\mathbf{Z}} \mathbf{C}
$$

to permit complex coefficients. $K(X ; \mathbf{C})$ is the free $\mathbf{C}$-vector space generated by $\operatorname{Vect}(X)$ subject to the relations $V \oplus W=V+W$. By using complex coefficients, we eliminate torsion which makes calculations much simpler. The Chern character is a morphism:

$$
c h: \operatorname{Vect}(X) \rightarrow H^{\text {even }}(X ; \mathbf{C})=\bigoplus_{q} H^{2 q}(X ; \mathbf{C})
$$

We define ch using characteristic classes in the second section if $X$ is a smooth manifold; it is possible to extend this definition using topological methods to more general topological settings. The identities:

$$
\operatorname{ch}(V \oplus W)=\operatorname{ch}(V)+\operatorname{ch}(W) \quad \text { and } \quad \operatorname{ch}(V \otimes W)=\operatorname{ch}(V) \operatorname{ch}(W)
$$

imply that we can extend:

$$
c h: K(X) \rightarrow H^{\text {even }}(X ; \mathbf{C})
$$

to be a ring homomorphism. We tensor this $\mathbf{Z}$-linear map with $\mathbf{C}$ to get

$$
c h: K(X ; \mathbf{C}) \rightarrow H^{\text {even }}(X ; \mathbf{C}) .
$$

Lemma 3.8.3 (Chern isomorphism). ch: $K(X ; \mathbf{C}) \rightarrow H^{\text {even }}(X ; \mathbf{C})$ is a ring isomorphism.

If $\widetilde{H}^{\text {even }}(X ; \mathbf{C})=\bigoplus_{q>o} H^{2 q}(X ; \mathbf{C})$ is the reduced even dimensional cohomology, then $c h: \widetilde{K}(X ; \mathbf{C}) \rightarrow \widetilde{H}^{\text {even }}(X ; \mathbf{C})$ is a ring isomorphism. For this reason, $\widetilde{K}(X)$ is often refered to as reduced $K$-theory. We emphasize that in this isomorphism we are ignoring torsion and that torsion is crucial
to understanding $K$-theory in general. Fortunately, the index is $\mathbf{Z}$-valued and we can ignore torsion in $K$-theory for an understanding of the index theorem.

We now return to studying the relation between $K(X)$ and $K(\Sigma X)$. We shall let $[X, Y]$ denote the set of homotopy classes of maps from $X$ to $Y$. We shall always assume that $X$ and $Y$ are equipped with base points and that all maps are basepoint preserving. We fix $2 k \geq \operatorname{dim} X$ and let $f: X \rightarrow S_{0}(2 k)$. Since $f(x)$ is self-adjoint, and $f(x)^{2}=I$, the eigenvalues of $f(x)$ are $\pm 1$. Since $\operatorname{Tr} f(x)=0$, each eigenvalue appears with multiplicity $k$. We let $\Pi_{ \pm}(f)$ be the bundles over $X$ which are sub-bundles of $X \times 1^{k}$ so that the fiber of $\Pi_{ \pm}(f)$ at $x$ is just the $\pm 1$ eigenspace of $f(x)$. If we define $\pi_{ \pm}(f)(x)=\frac{1}{2}(1 \pm f(x))$ then these are projections of constant rank $k$ with range $\Pi_{ \pm}(f)$. If $f$ and $f_{1}$ are homotopic maps, then they determine isomorphic vector bundles. Thus the assignment $f \mapsto \Pi_{+}(f) \in \operatorname{Vect}_{k}(X)$ defines a map $\Pi_{+}:\left[X, S_{0}(2 k)\right] \rightarrow \operatorname{Vect}_{k}(X)$.
Lemma 3.8.4. The natural map $\left[X, S_{0}(2 k)\right] \rightarrow \operatorname{Vect}_{k}(X)$ is bijective for $2 k \geq \operatorname{dim} X$.
Proof: Given $V \in \operatorname{Vect}_{k}(X)$ we choose $W \in \operatorname{Vect}_{k}(X)$ so $V \oplus W=1^{2 k}$. We choose fiber metrics on $V$ and on $W$ and make this sum orthogonal. By applying the Gram-Schmidt process to the given global frame, we can assume that there is an orthonormal global frame and consequently that $V \oplus W=1^{2 k}$ is an orthogonal direct sum. We let $\pi_{+}(x)$ be orthogonal projection on the fiber $V_{x}$ of $1^{2 k}$ and $f(x)=2 \pi_{+}(x)-I$. Then it is immediate that $f: X \rightarrow S_{0}(2 k)$ and $\pi_{+}(f)=V$. This proves the map is subjective. The injectivity comes from the same sorts of considerations as were used to prove Lemma 3.8.1 and is therefore omitted.

If $f: X \rightarrow \mathrm{GL}^{\prime}(k, \mathbf{C})$, we can extend the definition to let $\Pi_{ \pm}(f)(x)$ be the span of the generalized eigenvectors of $f$ corresponding to eigenvalues with positive/negative real part. Since there are no purely imaginary eigenvalues, $\Pi_{ \pm}(f)$ has constant rank and gives a vector bundle over $X$.

In a similar fashion, we can classify $\operatorname{Vect}_{k} \Sigma X=[X, \mathrm{U}(k)]$. Since $\mathrm{U}(k)$ is a deformation retract of $\mathrm{GL}(k, \mathbf{C})$, we identify $[X, \mathrm{U}(k)]=[X, \mathrm{GL}(k, \mathbf{C})]$. If $g: X \rightarrow \mathrm{GL}(k, \mathbf{C})$, we use $g$ as a clutching function to define a bundle over $\Sigma X$. Over $D_{ \pm}(X)$, we take the bundles $D_{ \pm}(X) \times \mathbf{C}^{k} . D_{+}(X) \cap D_{-}(X)=$ $X$. On the overlap, we identify $(x, z)_{+}=\left(x, z^{\prime}\right)_{-}$if $z \cdot g(x)=z^{\prime}$. If we let $s^{+}$and $s^{-}$be the usual frames for $C^{k}$ over $D_{ \pm}$then $\sum z_{i}^{+} g_{i j} s_{j}^{+}=\sum z_{i}^{-} s_{i}^{-}$ so that we identify the frames using the identity

$$
s^{-}=g s^{+}
$$

We denote this bundle by $V_{g}$. Homotopic maps define isomorphic bundles so we have a map $[X, \mathrm{GL}(k, \mathbf{C})] \rightarrow \operatorname{Vect}_{k} \Sigma X$. Conversely, given a vector bundle $V$ over $\Sigma X$ we can always choose local trivializations for $V$ over
$D_{ \pm}(X)$ since these spaces are contractible. The transition function $s^{-}=$ $g s^{+}$gives a map $g: X \rightarrow \mathrm{GL}(k, \mathbf{C})$. It is convenient to assume $X$ has a base point $x_{0}$ and to choose $s^{-}=s^{+}$at $x_{0}$. Thus $g\left(x_{0}\right)=I$. This shows the map $[X, \mathrm{GL}(k, \mathbf{C})] \rightarrow \operatorname{Vect}_{k} \Sigma X$ is surjective. If we had chosen different trivializations $\tilde{s}^{+}=h^{+} s^{+}$and $\tilde{s}^{-}=h^{-} s^{-}$then we would have obtained a new clutching function $\tilde{g}=h^{-} g\left(h^{+}\right)^{-1}$. Since $h^{ \pm}$are defined on contractible sets, they are null homotopic so $\tilde{g}$ is homotopic to $g$. This proves:

Lemma 3.8.5. The map $[X, \mathrm{GL}(k, \mathbf{C})]=[X, \mathrm{U}(k)] \rightarrow \operatorname{Vect}_{k}(\Sigma X)$, given by associating to a map $g$ the bundle defined by the clutching function $g$, is bijective.

It is always somewhat confusing to try to work directly with this definition. It is always a temptation to confuse the roles of $g$ and of its inverse as well as the transposes involved. There is another definition which avoids this difficulty and which will be very useful in computing specific examples. If $g: X \rightarrow \mathrm{GL}(k, \mathbf{C})$, we shall let $g(x) z$ denote matrix multiplication. We shall regard $\mathbf{C}^{k}$ as consisting of column vectors and let $g$ act as a matrix from the left. This is, of course, the opposite convention from that used previously.

Let $\theta \in\left[-\frac{\pi}{2}, \frac{\pi}{2}\right]$ be the suspension parameter. We define:

$$
\begin{aligned}
& \Sigma:[X, \mathrm{GL}(k, \mathbf{C})] \rightarrow\left[\Sigma X, \mathrm{GL}^{\prime}(2 k, \mathbf{C})\right] \\
& \Sigma:\left[X, \mathrm{GL}^{\prime}(k, \mathbf{C})\right] \rightarrow[\Sigma X, \mathrm{GL}(k, \mathbf{C})]
\end{aligned}
$$

by

$$
\begin{aligned}
\Sigma g(x, \theta) & =\left(\begin{array}{cc}
(\sin \theta) I_{k} & (\cos \theta) g^{*}(x) \\
(\cos \theta) g(x) & -(\sin \theta) I_{k}
\end{array}\right) \\
\Sigma f(x, \theta) & =(\cos \theta) f(x)-i(\sin \theta) I_{k}
\end{aligned}
$$

We check that $\Sigma$ has the desired ranges as follows. If $g: X \rightarrow \mathrm{GL}(k, \mathbf{C})$, then it is immediate that $\Sigma g$ is self-adjoint. We compute:

$$
(\Sigma g)^{2}=\left(\begin{array}{cc}
\left(\sin ^{2} \theta\right) I_{k}+\left(\cos ^{2} \theta\right) g^{*} g & 0 \\
0 & \left(\sin ^{2} \theta\right) I_{k}+\left(\cos ^{2} \theta\right) g g^{*}
\end{array}\right)
$$

This is non-singular since $g$ is invertible. Therefore $\Sigma g$ is invertible. If $\theta=\pi / 2$, then $\Sigma g=\left(\begin{array}{rr}I & 0 \\ 0 & -I\end{array}\right)$ is independent of $x$. If $g$ is unitary, $\Sigma g \in$ $\mathrm{S}_{0}(2 k)$. If $f \in \mathrm{GL}^{\prime}(k, \mathbf{C})$, then $f$ has no purely imaginary eigenvalues so $\Sigma f$ is non-singular.

Lemma 3.8.6. Let $g: X \rightarrow \mathrm{GL}(k, \mathbf{C})$ and construct the bundle $\Pi_{+}(\Sigma g)$ over $\Sigma X$. Then this bundle is defined by the clutching function $g$.

Proof: We may replace $g$ by a homotopic map without changing the isomorphism class of the bundle $\Pi_{+}(\Sigma g)$. Consequently, we may assume without loss of generality that $g: X \rightarrow \mathrm{U}(k)$ so $g^{*} g=g g^{*}=I_{k}$. Consequently $(\Sigma g)^{2}=I_{2 k}$ and $\pi_{+}(\Sigma g)=\frac{1}{2}\left(I_{2 k}+\Sigma g\right)$. If $z \in \mathbf{C}^{k}$, then:

$$
\pi_{+}(\Sigma g)(x, \theta)\binom{z}{0}=\frac{1}{2}\binom{z+(\sin \theta) z}{(\cos \theta) g(x) z}
$$

This projection from $\mathbf{C}^{k}$ to $\Pi_{+}(\Sigma g)$ is non-singular away from the south pole $S$ and can be used to give a trivialization of $\Pi_{+}(\Sigma g)$ on $\Sigma X-S$.

From this description, it is clear that $\Pi_{+}(\Sigma g)$ is spanned by vectors of the form $\frac{1}{2}\binom{(1+(\sin \theta)) z}{(\cos \theta) g(x) z}$ away from the south pole. At the south pole, $\Pi_{+}(\Sigma g)$ consists of all vectors of the form $\binom{0}{w}$. Consequently, projection on the second factor $\pi_{+}(S):\binom{a}{b} \rightarrow\binom{0}{b}$ is non-singular away from the north pole $N$ and gives a trivialization of $\Pi_{+}(\Sigma g)$ on $\Sigma X-N$. We restrict to the equator $X$ and compute the composite of these two maps to determine the clutching function:

$$
\binom{z}{0} \mapsto \frac{1}{2}\binom{z}{g(x) z} \mapsto \frac{1}{2}\binom{0}{g(x) z} .
$$

The function $g(x) / 2$ is homotopic to $g$ which completes the proof.
This is a very concrete description of the bundle defined by the clutching function $g$. In the examples we shall be considering, it will come equipped with a natural connection which will make computing characteristic classes much easier.

We now compute the double suspension. Fix $f: X \rightarrow \mathrm{~S}_{0}(2 k)$ and $g: X \rightarrow$ $\mathrm{GL}(k, \mathbf{C})$.

$$
\begin{aligned}
& \Sigma^{2} f(x, \theta, \phi) \\
& \quad=\left(\begin{array}{cc}
\sin \phi & \cos \phi\left\{(\cos \theta) f^{*}+i \sin \theta\right\} \\
\cos \phi\{(\cos \theta) f(x)-i \sin \theta\} & -\sin \phi
\end{array}\right) \\
& \Sigma^{2} g(x, \theta, \phi) \\
& \quad=\left(\begin{array}{cc}
\cos \phi \sin \theta-i \sin \phi & \cos \phi \cos \theta g^{*}(x) \\
\cos \phi \cos \theta g(x) & -\cos \phi \cos \theta-i \sin \phi
\end{array}\right)
\end{aligned}
$$

We let $\mathrm{U}(\infty)$ be the direct limit of the inclusions $\mathrm{U}(k) \rightarrow \mathrm{U}(k+1) \cdots$ and $\mathrm{S}_{0}(\infty)$ be the direct limit of the inclusions $\mathrm{S}_{0}(2 k) \rightarrow \mathrm{S}_{0}(2 k+2) \rightarrow \cdots$. Then we have identified:

$$
\widetilde{K}(X)=\left[X, \mathrm{~S}_{0}(\infty)\right] \quad \text { and } \quad \widetilde{K}(\Sigma X)=\left[\Sigma X, \mathrm{~S}_{0}(\infty)\right]=[S, \mathrm{U}(\infty)]
$$

We can now state:
Theorem 3.8.7 (Bott Periodicity). The map

$$
\Sigma^{2}: \widetilde{K}(X)=\left[X, \mathrm{~S}_{0}(\infty)\right] \rightarrow \widetilde{K}\left(\Sigma^{2} X\right)=\left[\Sigma^{2} X, \mathrm{~S}_{0}(\infty)\right]
$$

induces a ring isomorphism. Similarly, the map

$$
\Sigma^{2}: \widetilde{K}(\Sigma X)=[X, \mathrm{U}(\infty)] \rightarrow \widetilde{K}\left(\Sigma^{2} X\right)=\left[\Sigma^{2} X, \mathrm{U}(\infty)\right]
$$

is a ring isomorphism.
We note that $[X, \mathrm{U}(\infty)]$ inherits a natural additive structure from the group structure on $\mathrm{U}(\infty)$ by letting $g \oplus g^{\prime}$ be the direct sum of these two maps. This group structure is compatible with the additive structure on $\widetilde{K}(\Sigma X)$ since the clutching function of the direct sum is the direct sum of the clutching functions. Similarly, we can put a ring structure on $[X, \mathrm{U}(\infty)]$ using tensor products to be compatible with the ring structure on $\widetilde{K}(\Sigma X)$.

The Chern character identifies $\widetilde{K}(X ; \mathbf{C})$ with $\widetilde{H}^{\text {even }}(X ; \mathbf{C})$. We may identify $\Sigma^{2} X$ with a certain quotient of $X \times S^{2}$. Bott periodicity in this context becomes the assertion $K\left(X \times S^{2}\right)=K(X) \otimes K\left(S^{2}\right)$ which is the Kunneth formula in cohomology.

We now consider the case of a sphere $X=S^{n}$. The unitary group $\mathrm{U}(2)$ decomposes as $\mathrm{U}(2)=\mathrm{U}(1) \times \mathrm{SU}(2)=S^{1} \times S^{3}$ topologically so $\pi_{1}(\mathrm{U}(2))=$ $\mathbf{Z}, \pi_{2}(\mathrm{U}(2))=0$ and $\pi_{3}(\mathrm{U}(2))=\mathbf{Z}$. This implies $\widetilde{K}\left(S^{1}\right)=\widetilde{K}\left(S^{3}\right)=0$ and $\widetilde{K}\left(S^{2}\right)=\widetilde{K}\left(S^{4}\right)=\mathbf{Z}$. Using Bott periodicity, we know more generally that:

Lemma 3.8.8 (Bott periodicity). If $n$ is odd, then

$$
\widetilde{K}\left(S^{n}\right) \simeq \pi_{n-1}(\mathrm{U}(\infty))=0
$$

If $n$ is even, then

$$
\widetilde{K}\left(S^{n}\right) \simeq \pi_{n-1}(\mathrm{U}(\infty))=\mathbf{Z}
$$

It is useful to construct explicit generators of these groups. Let $x=$ $\left(x_{0}, \ldots, x_{n}\right) \in S^{n}$, let $y=\left(x, x_{n+1}\right) \in S^{n+1}$, and let $z=\left(y, x_{n+2}\right) \in$ $S^{n+2}$. If $f: S^{n} \rightarrow \mathrm{GL}^{\prime}(k, \mathbf{C})$ and $g: X \rightarrow \mathrm{GL}(k, \mathbf{C})$ we extend these to
be homogeneous of degree 1 with values in the $k \times k$ matrices. Then we compute:

$$
\begin{aligned}
\Sigma f(y) & =f(x)-i x_{n+1} \\
\Sigma^{2} f(z) & =\left(\begin{array}{cc}
x_{n+2} & f^{*}(x)+i x_{n+1} \\
f(x)-i x_{n+1} & -x_{n+2}
\end{array}\right) \\
\Sigma g(y) & =\left(\begin{array}{cc}
x_{n+1} & g^{*}(x) \\
g(x) & -x_{n+1}
\end{array}\right) \\
\Sigma^{2} g(z) & =\left(\begin{array}{cc}
x_{n+1}-i x_{n+2} & g^{*}(x) \\
g(x) & -x_{n+1}-i x_{n+2}
\end{array}\right) .
\end{aligned}
$$

If we suppose that $f$ is self-adjoint, then we can express:

$$
\Sigma^{2} f(z)=x_{n+2}\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right) \otimes I_{k}+x_{n+1}\left(\begin{array}{rr}
0 & i \\
-i & 0
\end{array}\right) \otimes I_{k}+\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right) \otimes f(x)
$$

We can now construct generators for $\pi_{n-1} \mathrm{U}(\infty)$ and $\widetilde{K}\left(S^{n}\right)$ using Clifford algebras. Let $g\left(x_{0}, x_{1}\right)=x_{0}-i x_{1}$ generate $\pi_{1}\left(S^{1}\right)=\mathbf{Z}$ then

$$
\Sigma g\left(x_{0}, x_{1}, x_{2}\right)=x_{0} e_{0}+x_{1} e_{1}+x_{2} e_{2}
$$

for

$$
e_{0}=\left(\begin{array}{rr}
0 & 1 \\
1 & 0
\end{array}\right), \quad e_{1}=\left(\begin{array}{rr}
0 & i \\
-i & 0
\end{array}\right), \quad e_{2}=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right)
$$

The $\left\{e_{j}\right\}$ satisfy the relations $e_{j} e_{k}+e_{k} e_{j}=2 \delta_{j k}$ and $e_{0} e_{1} e_{2}=-i I_{2} . \Pi_{+}(\Sigma g)$ is a line bundle over $S^{2}$ which generates $\widetilde{K}\left(S^{2}\right)$. We compute:

$$
\Sigma^{2} g\left(x_{0}, x_{1}, x_{2}, x_{3}\right)=x_{0} e_{0}+x_{1} e_{1}+x_{2} e_{2}-i x_{3} I_{2}
$$

If we introduce $z_{1}=x_{0}+i x_{1}$ and $z_{2}=x_{2}+i x_{3}$ then

$$
\Sigma^{2} g\left(z_{1}, z_{2}\right)=\left(\begin{array}{rr}
\bar{z}_{2} & z_{1} \\
\bar{z}_{1} & -z_{2}
\end{array}\right)
$$

Consequently $\Sigma^{2} g$ generates $\pi_{3}(\mathrm{U}(\infty))$. We suspend once to construct:

$$
\begin{aligned}
& \Sigma^{3} g\left(x_{0}, x_{1}, x_{2}, x_{3}, x_{4}\right) \\
& \quad=x_{0} e_{0} \otimes e_{0}+x_{1} e_{0} \otimes e_{1}+x_{2} e_{0} \otimes e_{2}+x_{3} e_{1} \otimes I+x_{4} e_{2} \otimes I
\end{aligned}
$$

The bundle $\Pi_{+}\left(\Sigma^{3} g\right)$ is a 2-plane bundle over $S^{4}$ which generates $\widetilde{K}\left(S^{4}\right)$. We express $\Sigma^{3} g=x_{0} e_{0}^{4}+\cdots+x_{4} e_{4}^{4}$ then these matrices satisfy the commutation relations $e_{j}^{4} e_{j}^{4}+e_{k}^{4} e_{j}^{4}=2 \delta_{j k}$ and $e_{0}^{4} e_{1}^{4} e_{2}^{4} e_{3}^{4} e_{4}^{4}=-I$.

We proceed inductively to define matrices $e_{j}^{2 k}$ for $0 \leq j \leq 2 k$ so that $e_{i} e_{j}+e_{j} e_{i}=2 \delta_{i j}$ and $e_{0}^{2 k} \ldots e_{2 k}^{2 k}=(-i)^{k} I$. These are matrics of shape $2^{k} \times 2^{k}$ such that $\Sigma^{2 k-1}(g)(x)=\sum x_{j} e_{j}^{2 k}$. In Lemma 2.1.5 we computed that:

$$
\int_{S^{2 k}} c h_{k}\left(\Pi_{+} \Sigma^{2 k-1} g\right)=i^{k} 2^{-k} \operatorname{Tr}\left(e_{0}^{2 k} \ldots e_{2 k}^{2 k}\right)=1
$$

These bundles generate $\widetilde{K}\left(S^{2 k}\right)$ and $\Sigma^{2 k}(g)$ generates $\pi_{2 k+1}(\mathrm{U}(\infty))$. We summarize these calculations as follows:

Lemma 3.8.9. Let $\left\{e_{0}, \ldots, e_{2 k}\right\}$ be a collection of self-adjoint matrices of shape $2^{k} \times 2^{k}$ such that $e_{i} e_{j}+e_{j} e_{i}=2 \delta_{i j}$ and such that $e_{0} \ldots e_{2 k}=(-i)^{k} I$. We define $e(x)=x_{0} e_{0}+\cdots+x_{2 k} e_{2 k}$ for $x \in S^{2 k}$. Let $\Pi_{+}(e)$ be the bundle of +1 eigenvectors of e over $S^{2 k}$. Then $\Pi_{+}(e)$ generates $\widetilde{K}\left(S^{2 k}\right) \simeq \mathbf{Z} . \Sigma e(y)=$ $e(x)-i x_{2 k+1}$ generates $\pi_{2 k+1}(\mathrm{U}(\infty)) \simeq \mathbf{Z} . \int_{S^{2 k}} c h_{k}(V): \widetilde{K}\left(S^{2 k}\right) \rightarrow \mathbf{Z}$ is an isomorphism and $\int_{S^{2 k}} c h_{k}\left(\Pi_{+} e\right)=1$.
Proof: We note that $\int_{S^{2 k}} c h_{k}(V)$ is the index of the spin complex with coefficients in $V$ since all the Pontrjagin forms of $T\left(S^{2 k}\right)$ vanish except for $p_{0}$. Thus this integral, called the topological charge, is always an integer. We have checked that the integral is 1 on a generator and hence the map is surjective. Since $\widetilde{K}\left(S^{2 k}\right)=\mathbf{Z}$, it must be bijective.

It was extremely convenient to have the bundle with clutching function $\Sigma^{2 k-2} g$ so concretely given so that we could apply Lemma 2.1.5 to compute the topological charge. This will also be important in the next chapter.

### 3.9. The Atiyah-Singer Index Theorem.

In this section, we shall discuss the Atiyah-Singer theorem for a general elliptic complex by interpreting the index as a map in $K$-theory. Let $M$ be smooth, compact, and without boundary. For the moment we make no assumptions regarding the parity of the dimension $m$. We do not assume $M$ is orientable. Let $P: C^{\infty}\left(V_{1}\right) \rightarrow \mathbf{C}^{\infty}\left(V_{2}\right)$ be an elliptic complex with leading symbol $p(x, \xi): S\left(T^{*} M\right) \rightarrow \operatorname{HOM}\left(V_{1}, V_{2}\right)$. We let $\Sigma\left(T^{*} M\right)$ be the fiberwise suspension of the unit sphere bundle $S\left(T^{*} M\right)$. We identify $\Sigma\left(T^{*} M\right)=S\left(T^{*} M \oplus 1\right)$. We generalize the construction of section 3.8 to define $\Sigma p: \Sigma\left(T^{*} M\right) \rightarrow \operatorname{END}\left(V_{1} \oplus V_{2}\right)$ by:

$$
\Sigma p(x, \xi, \theta)=\left(\begin{array}{cc}
(\sin \theta) I_{V_{1}} & (\cos \theta) p^{*}(x, \xi) \\
(\cos \theta) p(x, \xi) & -(\sin \theta) I_{V_{2}}
\end{array}\right)
$$

This is a self-adjoint invertible endomorphism. We let $\Pi_{ \pm}(\Sigma p)$ be the subbundle of $V_{1} \oplus V_{2}$ over $\Sigma\left(T^{*} M\right)$ corresponding to the span of the eigenvectors of $\Sigma p$ with positive/negative eigenvalues. If we have given connections on $V_{1}$ and $V_{2}$, we can project these connections to define natural connections on $\Pi_{ \pm}(\Sigma p)$. The clutching function of $\Pi_{+}(\Sigma p)$ is $p$ in a sense we explain as follows:

We form the disk bundles $D_{ \pm}(M)$ over $M$ corresponding to the northern and southern hemispheres of the fiber spheres of $\Sigma\left(T^{*} M\right)$. Lemma 3.8.6 generalizes immediately to let us identify $\Pi_{+}(\Sigma p)$ with the bundle $V_{1}^{+} \cup V_{2}^{-}$over the disjoint union $D_{+}(M) \cup D_{-}(M)$ attached using the clutching function $p$ over their common boundary $S\left(T^{*} M\right)$. If $\operatorname{dim} V_{1}=k$, then $\Pi_{+}(\Sigma p) \in \operatorname{Vect}_{k}\left(\Sigma\left(T^{*} M\right)\right)$. Conversely, we suppose given a bundle $V \in \operatorname{Vect}_{k}\left(\Sigma\left(T^{*} M\right)\right)$. Let $N: M \rightarrow \Sigma\left(T^{*} M\right)$ and $S: M \rightarrow \Sigma\left(T^{*} M\right)$ be the natural sections mapping $M$ to the northern and southern poles of the fiber spheres; $N(x)=(x, 0,1)$ and $S(x)=(x, 0,-1)$ in $S\left(T^{*} M \oplus 1\right) . N$ and $S$ are the centers of the disk bundles $D_{ \pm}(M)$. We let $N^{*}(V)=V_{1}$ and $S^{*}(V)=V_{2}$ be the induced vector bundles over $M . D_{ \pm}(M)$ deformation retracts to $M \times\{N\}$ and $M \times\{S\}$. Thus $V$ restricted to $D_{ \pm}(M)$ is cannonically isomorphic to the pull back of $V_{1}$ and $V_{2}$. On the intersection $S\left(T^{*} M\right)=D_{+}(M) \cap D_{-}(M)$ we have a clutching or glueing function relating the two decompositions of $V$. This gives rise to a map $p: S\left(T^{*} M\right) \rightarrow \operatorname{HOM}\left(V_{1}, V_{2}\right)$ which is non-singular. The same argument as that given in the proof of Lemma 3.8.5 shows that $V$ is completely determined by the isomorphism class of $V_{1}$ and of $V_{2}$ together with the homotopy class of the map $p: S\left(T^{*} M\right) \rightarrow \operatorname{HOM}\left(V_{1}, V_{2}\right)$.

Given an order $\nu$ we can recover the leading symbol $p$ by extending $p$ from $S\left(T^{*} M\right)$ to $T^{*} M$ to be homogeneous of order $\nu$. We use this to define an elliptic operator $P_{\nu}: C^{\infty}\left(V_{1}\right) \rightarrow C^{\infty}\left(V_{2}\right)$ with leading symbol $p_{\nu}$. If $Q_{\nu}$ is another operator with the same leading symbol, then we define
$P_{\nu}(t)=t P_{\nu}+(1-t) Q_{\nu}$. This is a 1-parameter family of such operators with the same leading symbol. Consequently by Lemma 1.4.4 index $\left(P_{\nu}\right)=$ index $\left(Q_{\nu}\right)$. Similarly, if we replace $p$ by a homotopic symbol, then the index is unchanged. Finally, suppose we give two orders of homgeneity $\nu_{1}>\nu_{2}$. We can choose a self-adjoint pseudo-differential operator $R$ on $C^{\infty}\left(V_{2}\right)$ with leading symbol $|\xi|^{\nu_{1}-\nu_{2}} I_{V_{2}}$. Then we can let $P_{\nu_{1}}=R P_{\nu_{2}}$ so $\operatorname{index}\left(P_{\nu_{1}}\right)=\operatorname{index}(R)+\operatorname{index}\left(P_{\nu_{2}}\right)$. Since $R$ is self-adjoint, its index is zero so index $\left(P_{\nu_{1}}\right)=\operatorname{index}\left(P_{\nu_{2}}\right)$. This shows that the index only depends on the homotopy class of the clutching map over $S\left(T^{*} M\right)$ and is independent of the order of homogeneity and the extension and the particular operator chosen. Consequently, we can regard index: $\operatorname{Vect}\left(\Sigma\left(T^{*} M\right)\right) \rightarrow \mathbf{Z}$ so that $\operatorname{index}\left(\Pi_{+}(\Sigma p)\right)=\operatorname{index}(P)$ if $P$ is an elliptic operator. (We can always "roll up" an elliptic complex to give a 2 -term elliptic complex in computing the index so it suffices to consider this case).

It is clear from our definition that $\Sigma(p \oplus q)=\Sigma(p) \oplus \Sigma(q)$ and therefore $\Pi_{+}(\Sigma(p \oplus q))=\Pi_{+}(\Sigma(p)) \oplus \Pi_{+}(\Sigma(q))$. Since index $(P \oplus Q)=\operatorname{index}(P)+$ index $(Q)$ we conclude:

$$
\operatorname{index}(V \oplus W)=\operatorname{index}(V)+\operatorname{index}(W) \quad \text { for } V, W \in \operatorname{Vect}\left(\Sigma\left(T^{*} M\right)\right)
$$

This permits us to extend index: $K\left(\Sigma\left(T^{*} M\right)\right) \rightarrow \mathbf{Z}$ to be Z-linear. We tensor with the complex numbers to extend index: $K\left(\Sigma\left(T^{*} M\right) ; \mathbf{C}\right) \rightarrow \mathbf{C}$, so that:

Lemma 3.9.1. There is a natural map index: $K\left(\Sigma\left(T^{*} M\right) ; \mathbf{C}\right) \rightarrow \mathbf{C}$ which is linear so that $\operatorname{index}(P)=\operatorname{index}\left(\Pi_{+}(\Sigma p)\right)$ if $P: C^{\infty} V_{1} \rightarrow C^{\infty} V_{2}$ is an elliptic complex over $M$ with symbol $p$.

There is a natural projection map $\pi: \Sigma\left(T^{*} M\right) \rightarrow M$. This gives a natural map $\pi^{*}: K(M ; \mathbf{C}) \rightarrow K\left(\Sigma\left(T^{*} M\right) ; \mathbf{C}\right)$. If $N$ denotes the north pole section, then $\pi N=1_{M}$ so $N^{*} \pi^{*}=1$ and consequently $\pi^{*}$ is injective. This permits us to regard $K(M ; \mathbf{C})$ as a subspace of $K\left(\Sigma\left(T^{*} M\right) ; \mathbf{C}\right)$. If $V=\pi^{*} V_{1}$, then the clutching function defining $V$ is just the identity map. Consequently, the corresponding elliptic operator $P$ can be taken to be a self-adjoint operator on $C^{\infty}(V)$ which has index zero. This proves:

Lemma 3.9.2. If $V \in K\left(\Sigma\left(T^{*} M\right) ; \mathbf{C}\right)$ can be written as $\pi^{*} V_{1}$ for $V_{1} \in$ $K(M ; \mathbf{C})$ then index $(V)=0$. Thus index: $K\left(\Sigma\left(T^{*} M\right) ; \mathbf{C}\right) / K(M ; \mathbf{C}) \rightarrow \mathbf{C}$.

These two lemmas show that all the information contained in an elliptic complex from the point of view of computing its index is contained in the corresponding description in $K$-theory. The Chern character gives an isomorphism of $K(X ; \mathbf{C})$ to the even dimensional cohomology. We will exploit this isomorphism to give a formula for the index in terms of cohomology.

In addition to the additive structure on $K(X ; \mathbf{C})$, there is also a ring structure. This ring structure also has its analogue with respect to elliptic
operators as we have discussed previously. The multiplicative nature of the four classical elliptic complexes played a fundamental role in determining the normalizing constants in the formula for their index.

We give $\Sigma\left(T^{*} M\right)$ the simplectic orientation. If $x=\left(x_{1}, \ldots, x_{m}\right)$ is a system of local coordinates on $M$, let $\xi=\left(\xi_{1}, \ldots, \xi_{m}\right)$ be the fiber coordinates on $T^{*} M$. Let $u$ be the additional fiber coordinate on $T^{*} M \oplus 1$. We orient $T^{*} M \oplus 1$ using the form:

$$
\omega_{2 m+1}=d x_{1} \wedge d \xi_{1} \wedge \cdots \wedge d x_{m} \wedge d \xi_{m} \wedge d u
$$

Let $\vec{N}$ be the outward normal on $S\left(T^{*} M \oplus 1\right)$ and define $\omega_{2 m}$ on $S\left(T^{*} M \oplus 1\right)$ by:

$$
\omega_{2 m+1}=\vec{N} \wedge \omega_{2 m}=\omega_{2 m} \wedge \vec{N}
$$

This gives the orientation of Stokes theorem.
Lemma 3.9.3.
(a) Let $P: C^{\infty}\left(V_{1}\right) \rightarrow C^{\infty}\left(V_{2}\right)$ be an elliptic complex over $M_{1}$ and let $Q: C^{\infty}\left(W_{1}\right) \rightarrow C^{\infty}\left(W_{2}\right)$ be an elliptic complex over $M_{2}$. We assume $P$ and $Q$ are partial differential operators of the same order and we let $M=$ $M_{1} \times M_{2}$ and define over $M$

$$
\begin{aligned}
R=(P \otimes 1+1 \otimes Q) \oplus\left(P^{*} \otimes 1-1 \otimes Q^{*}\right) & : C^{\infty}\left(V_{1} \otimes W_{1} \oplus V_{2} \otimes W_{2}\right) \\
& \rightarrow C^{\infty}\left(V_{2} \otimes W_{1} \oplus V_{1} \otimes W_{2}\right) .
\end{aligned}
$$

Then $R$ is elliptic and $\operatorname{index}(R)=\operatorname{index}(P) \cdot \operatorname{index}(Q)$.
(b) The four classical elliptic complexes discussed earlier can be decomposed in this fashion over product manifolds.
(c) Let $p$ and $q$ be arbitrary elliptic symbols over $M_{1}$ and $M_{2}$ and define

$$
r=\left(\begin{array}{cc}
p \otimes 1 & -1 \otimes q^{*} \\
1 \otimes q & p^{*} \otimes 1
\end{array}\right) \quad \text { over } M=M_{1} \times M_{2}
$$

Let $\theta_{i} \in H^{\text {even }}\left(M_{i} ; \mathbf{C}\right)$ for $i=1,2$, then:

$$
\begin{aligned}
\int_{\Sigma\left(T^{*} M\right)} \theta_{1} & \wedge \theta_{2} \wedge \operatorname{ch}\left(\Pi_{+}(\Sigma r)\right) \\
& =\int_{\Sigma\left(T^{*} M_{1}\right)} \theta_{1} \wedge \operatorname{ch}\left(\Pi_{+}(\Sigma p)\right) \cdot \int_{\Sigma\left(T^{*} M_{2}\right)} \theta_{2} \wedge \operatorname{ch}\left(\Pi_{+}(\Sigma q)\right)
\end{aligned}
$$

(d) Let $q$ be an elliptic symbol over $M_{1}$ and let $p$ be a self-adjoint elliptic symbol over $M_{2}$. Over $M=M_{1} \times M_{2}$ define the self-adjoint elliptic symbol $r$ by:

$$
r=\left(\begin{array}{cc}
1 \otimes p & q^{*} \otimes 1 \\
q \otimes 1 & -1 \otimes p
\end{array}\right)
$$

Let $\theta_{1} \in \Lambda^{e}\left(M_{1}\right)$ and $\theta_{2} \in \Lambda^{\circ}\left(M_{2}\right)$ be closed differential forms. Give $S\left(T^{*} M\right), \Sigma\left(T^{*} M_{1}\right)$ and $S\left(T^{*} M_{2}\right)$ the orientations induced from the simplectic orientations. Then:

$$
\begin{aligned}
\int_{S\left(T^{*} M\right)} \theta_{1} \wedge & \theta_{2} \wedge c h\left(\Pi_{+} r\right) \\
& =\int_{\Sigma\left(T^{*} M_{1}\right)} \theta_{1} \wedge \operatorname{ch}\left(\Pi_{+}(\Sigma q)\right) \cdot \int_{S\left(T^{*} M_{2}\right)} \theta_{1} \wedge c h\left(\Pi_{+}(p)\right)
\end{aligned}
$$

Remark: We will use (d) to discuss the eta invariant in Chapter 4; we include this integral at this point since the proof is similar to that of (c).

Proof: We let ( $p, q, r$ ) be the symbols of the operators involved. Then:

$$
r=\left(\begin{array}{cc}
p \otimes 1 & -1 \otimes q^{*} \\
1 \otimes q & p^{*} \otimes 1
\end{array}\right) \quad r^{*}=\left(\begin{array}{cc}
p^{*} \otimes 1 & 1 \otimes q^{*} \\
-1 \otimes q & p \otimes 1
\end{array}\right)
$$

so that:

$$
\begin{aligned}
& r^{*} r=\left(\begin{array}{cc}
p^{*} p \otimes 1+1 \otimes q^{*} q & 0 \\
0 & p p^{*} \otimes 1+1 \otimes q q^{*}
\end{array}\right) \\
& r r^{*}=\left(\begin{array}{cc}
p p^{*} \otimes 1+1 \otimes q^{*} q & 0 \\
0 & p^{*} p \otimes 1+1 \otimes q q^{*}
\end{array}\right)
\end{aligned}
$$

$r^{*} r$ and $r r^{*}$ are positive self-adjoint matrices if $\left(\xi^{1}, \xi^{2}\right) \neq(0,0)$. This verifies the ellipticity. We note that if $(P, Q)$ are pseudo-differential, $R$ will still be formally elliptic, but the symbol will not in general be smooth at $\xi^{1}=0$ or $\xi^{2}=0$ and hence $R$ will not be a pseudo-differential operator in that case. We compute:

$$
\begin{aligned}
& R^{*} R=\left(\begin{array}{cc}
P^{*} P \otimes 1+1 \otimes Q^{*} Q & 0 \\
0 & P P^{*} \otimes 1+1 \otimes Q Q^{*}
\end{array}\right) \\
& R R^{*}=\left(\begin{array}{cc}
P P^{*} \otimes 1+1 \otimes Q^{*} Q & 0 \\
0 & P^{*} P \otimes 1+1 \otimes Q Q^{*}
\end{array}\right) \\
& \mathrm{N}\left(R^{*} R\right)= \mathrm{N}\left(P^{*} P\right) \otimes \mathrm{N}\left(Q^{*} Q\right) \oplus \mathrm{N}\left(P P^{*}\right) \otimes \mathrm{N}\left(Q Q^{*}\right) \\
& \mathrm{N}\left(R R^{*}\right)= \mathrm{N}\left(P P^{*}\right) \otimes \mathrm{N}\left(Q^{*} Q\right) \oplus \mathrm{N}\left(P^{*} P\right) \otimes \mathrm{N}\left(Q Q^{*}\right) \\
& \operatorname{index}(R)=\left\{\operatorname{dim\mathrm {N}}\left(P^{*} P\right)-\operatorname{dim} \mathrm{N}\left(P P^{*}\right)\right\} \\
& \times\left\{\operatorname{dim} \mathrm{N}\left(Q^{*} Q\right)-\operatorname{dim~} \mathrm{N}\left(Q Q^{*}\right)\right\} \\
&= \operatorname{index}(P) \operatorname{index}(Q)
\end{aligned}
$$

which completes the proof of (a).

We verify (b) on the symbol level. First consider the de Rham complex and decompose:

$$
\begin{aligned}
\Lambda\left(T^{*} M_{1}\right) & =\Lambda_{1}^{\mathrm{e}} \oplus \Lambda_{1}^{\mathrm{o}}, \quad \Lambda\left(T^{*} M_{2}\right)=\Lambda_{2}^{\mathrm{e}} \oplus \Lambda_{2}^{\mathrm{o}} \\
\Lambda\left(T^{*} M\right) & =\left(\Lambda_{1}^{\mathrm{e}} \otimes \Lambda_{2}^{\mathrm{e}} \oplus \Lambda_{1}^{\mathrm{o}} \otimes \Lambda_{2}^{\mathrm{o}}\right) \oplus\left(\Lambda_{1}^{\mathrm{o}} \otimes \Lambda_{2}^{\mathrm{e}} \oplus \Lambda_{1}^{\mathrm{e}} \otimes \Lambda_{2}^{\mathrm{o}}\right)
\end{aligned}
$$

Under this decomposition:

$$
\sigma_{L}\left((d+\delta)_{M}\right)\left(\xi^{1}, \xi^{2}\right)=\left(\begin{array}{cc}
c\left(\xi^{1}\right) \otimes 1 & -1 \otimes c\left(\xi^{2}\right) \\
1 \otimes c\left(\xi^{2}\right) & c\left(\xi^{1}\right) \otimes 1
\end{array}\right)
$$

$c(\cdot)$ denotes Clifford multiplication. This verifies the de Rham complex decomposes properly.

The signature complex is more complicated. We decompose

$$
\Lambda^{ \pm}\left(T^{*} M_{1}\right)=\Lambda_{1}^{ \pm, \mathrm{e}} \oplus \Lambda_{1}^{ \pm, \mathrm{o}}
$$

to decompose the signature complex into two complexes:

$$
\begin{aligned}
& (d+\delta): C^{\infty}\left(\Lambda_{1}^{+, \mathrm{e}}\right) \rightarrow C^{\infty}\left(\Lambda_{1}^{-, \mathrm{o}}\right) \\
& (d+\delta): C^{\infty}\left(\Lambda_{1}^{+, \mathrm{o}}\right) \rightarrow C^{\infty}\left(\Lambda_{1}^{-, \mathrm{e}}\right)
\end{aligned}
$$

Under this decomposition, the signature complex of $M$ decomposes into four complexes. If, for example, we consider the complex:

$$
\begin{aligned}
(d+\delta): C^{\infty}\left(\Lambda_{1}^{+, \mathrm{e}} \otimes \Lambda_{2}^{+, \mathrm{e}} \oplus \Lambda_{1}^{-, \mathrm{o}} \otimes\right. & \left.\Lambda_{2}^{-, \mathrm{o}}\right) \rightarrow \\
& C^{\infty}\left(\Lambda_{1}^{-, \mathrm{o}} \otimes \Lambda_{2}^{+, \mathrm{e}} \oplus \Lambda_{1}^{+, \mathrm{e}} \otimes \Lambda_{2}^{-, \mathrm{o}}\right)
\end{aligned}
$$

then the same argument as that given for the de Rham complex applies to show the symbol is:

$$
\left(\begin{array}{cc}
c\left(\xi^{1}\right) \otimes 1 & -1 \otimes c\left(\xi^{2}\right) \\
1 \otimes c\left(\xi^{2}\right) & c\left(\xi^{1}\right) \otimes 1
\end{array}\right)
$$

If we consider the complex:

$$
\begin{aligned}
(d+\delta): C^{\infty}\left(\Lambda_{1}^{+, \mathrm{o}} \otimes \Lambda_{2}^{+\mathrm{e}} \oplus \Lambda_{1}^{-, \mathrm{e}}\right. & \left.\otimes \Lambda_{2}^{-, \mathrm{o}}\right) \\
& \rightarrow C^{\infty}\left(\Lambda_{1}^{-, \mathrm{e}} \otimes \Lambda_{2}^{+, \mathrm{e}} \oplus \Lambda_{1}^{+, \mathrm{o}} \otimes \Lambda_{2}^{-, \mathrm{o}}\right)
\end{aligned}
$$

then we conclude the symbol is

$$
\left(\begin{array}{cc}
c\left(\xi^{1}\right) \otimes 1 & 1 \otimes c\left(\xi^{2}\right) \\
-1 \otimes\left(\xi^{2}\right) & c\left(\xi^{1}\right) \otimes 1
\end{array}\right)
$$

which isn't right. We can adjust the sign problem for $\xi^{2}$ either by changing one of the identifications with $\Lambda^{+}(M)$ or by changing the $\operatorname{sign}$ of $Q$ (which won't affect the index). The remaining two cases are similar. The spin and Dolbeault complexes are also similar. If we take coefficients in an auxilary bundle $V$, the symbols involved are unchanged and the same arguments hold. This proves (b).

The proof of (c) is more complicated. We assume without loss of generality that $p$ and $q$ are homogeneous of degree 1 . Let $\left(\xi^{1}, \xi^{2}, u\right)$ parametrize the fibers of $T^{*} M \oplus 1$. At $\xi^{1}=(1 / \sqrt{2}, 0, \ldots, 0), \xi^{2}=(1 / \sqrt{2}, 0, \ldots, 0)$, $u=0$ the orientation is given by:

$$
\begin{aligned}
-d x_{1}^{1} \wedge d \xi_{1}^{1} & \wedge d x_{2}^{1} \wedge d \xi_{2}^{1} \wedge \cdots \\
& \wedge d x_{m_{1}}^{1} \wedge d \xi_{m_{1}}^{1} \wedge d x_{1}^{2} \wedge d x_{2}^{2} \wedge d \xi_{2}^{2} \wedge \cdots \wedge d x_{m_{2}}^{2} \wedge d \xi_{m_{2}}^{2} \wedge d u
\end{aligned}
$$

We have omitted $d \xi_{2}^{1}$ (and changed the sign) since it points outward at this point of the sphere to get the orientation on $\Sigma\left(T^{*} M\right)$.

In matrix form we have on $\left(V_{1} \otimes W_{1}\right) \oplus\left(V_{2} \otimes W_{2}\right) \oplus\left(V_{2} \otimes W_{1}\right) \oplus\left(V_{1} \otimes W_{2}\right)$ that:

$$
\Sigma r=\left(\begin{array}{cc}
u\left(\begin{array}{cc}
1 & 0 \\
0 & 1
\end{array}\right) & \left(\begin{array}{cc}
p^{*} \otimes 1 & 1 \otimes q^{*} \\
-1 \otimes q & p \otimes 1
\end{array}\right) \\
\left(\begin{array}{cc}
p \otimes 1 & -1 \otimes q^{*} \\
1 \otimes q & p^{*} \otimes 1
\end{array}\right) & -u\left(\begin{array}{cc}
1 & 0 \\
0 & 1
\end{array}\right)
\end{array}\right) .
$$

This is not a very convenient form to work with. We define

$$
\begin{aligned}
& \gamma_{1}=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right) \quad \text { and } \quad \alpha=\left(\begin{array}{cc}
0 & p^{*} \\
p & 0
\end{array}\right) \quad \text { on } \quad V_{1} \oplus V_{2}=V \\
& \gamma_{2}=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right) \quad \text { and } \quad \beta=\left(\begin{array}{cc}
0 & q^{*} \\
q & 0
\end{array}\right) \quad \text { on } W_{1} \oplus W_{2}=W
\end{aligned}
$$

and compute that $\Sigma r=u \gamma_{1} \otimes \gamma_{2}+\alpha \otimes I_{W}+\gamma_{1} \otimes \beta$. We replace $p$ and $q$ by homotopic symbols so that:

$$
p^{*} p=\left|\xi^{1}\right|^{2} I_{V_{1}}, \quad p p^{*}=\left|\xi^{1}\right|^{2} I_{V_{2}}, \quad q^{*} q=\left|\xi^{2}\right|^{2} I_{W_{1}}, \quad q q^{*}=\left|\xi^{2}\right|^{2} I_{W_{2}}
$$

Since $\left\{\gamma_{1} \otimes \gamma_{2}, \alpha \otimes I_{W}, \gamma_{1} \otimes \beta\right\}$ all anti-commute and are self-adjoint,

$$
(\Sigma r)^{2}=\left(u^{2}+\left|\xi^{1}\right|^{2}+\left|\xi^{2}\right|^{2}\right) I \quad \text { on } V \otimes W
$$

We parametrize $\Sigma\left(T^{*} M\right)$ by $S\left(T^{*} M_{1}\right) \times\left[0, \frac{\pi}{2}\right] \times \Sigma\left(T^{*} M_{2}\right)$ in the form:

$$
\left(\xi^{1} \cos \theta, \xi^{2} \sin \theta, u \sin \theta\right)
$$

We compute the orientation by studying $\theta=\frac{\pi}{4}, \xi^{1}=(1,0, \ldots, 0), \xi^{2}=$ $(1,0, \ldots, 0), u=0$. Since $d \xi_{1}^{1}=-\sin \theta d \theta$, the orientation is given by:

$$
\begin{aligned}
d x_{1}^{1} \wedge d \theta & \wedge d x_{2}^{1} \wedge d \xi_{2}^{1} \wedge \cdots \\
& \wedge d x_{m_{1}}^{1} \wedge d \xi_{m_{1}}^{1} \wedge d x_{1}^{2} \wedge d x_{2}^{2} \wedge d \xi_{2}^{2} \wedge \cdots \wedge d x_{m_{2}}^{2} \wedge d \xi_{m_{2}}^{2} \wedge d u
\end{aligned}
$$

If we identify $S\left(T^{*} M_{1}\right) \times\left[0, \frac{\pi}{2}\right]$ with $D_{+}\left(T^{*} M_{1}\right)$ then this orientation is:

$$
\omega_{2 m_{1}}^{1} \wedge \omega_{2 m_{2}}^{2}
$$

so the orientations are compatible.
In this parametrization, we compute $\Sigma r\left(\xi^{1}, \theta, \xi^{2}, u\right)=(\sin \theta) \gamma_{1} \otimes\left(u \gamma_{2}+\right.$ $\left.\beta\left(\xi^{2}\right)\right)+(\cos \theta) \alpha\left(\xi^{1}\right) \otimes I_{W}$. Since $\Sigma q=u \gamma_{2}+\beta\left(\xi^{2}\right)$ satisfies $(\Sigma q)^{2}=$ $\left|\xi^{2}\right|^{2}+u_{2}^{2}$ on $W$, we may decompose $W=\Pi_{+}(\Sigma q) \oplus \Pi_{-}(\Sigma q)$. Then:

$$
\begin{array}{ll}
\Sigma r=\left\{(\sin \theta) \gamma_{1}+(\cos \theta) \alpha\left(\xi^{1}\right)\right\} \otimes I=\Sigma p\left(\xi^{1}, \theta\right) \otimes I & \text { on } V \otimes \Pi_{+}(\Sigma q) \\
\Sigma r=\left\{\sin (-\theta) \gamma_{1}+(\cos \theta) \alpha\left(\xi^{1}\right)\right\} \otimes I=\Sigma p\left(\xi^{1},-\theta\right) \otimes I & \text { on } V \otimes \Pi_{-}(\Sigma q)
\end{array}
$$

Consequently:

$$
\begin{aligned}
& \Pi_{+}(\Sigma r)=\left\{\Pi_{+}(\Sigma p)\left(\xi^{1}, \theta\right) \otimes \Pi_{+}(\Sigma q)\right\} \oplus\left\{\Pi_{+}(\Sigma p)\left(\xi^{1},-\theta\right) \otimes \Pi_{-}(\Sigma q)\right\} \\
& \text { over }\left(D_{+} M_{1}\right) \times \Sigma\left(T^{*} M_{2}\right)
\end{aligned}
$$

If we replace $-\theta$ by $\theta$ in the second factor, we may replace $\Pi_{+}(\Sigma p)\left(\xi^{1},-\theta\right) \otimes$ $\Pi_{-}(\Sigma q)$ by $\Pi_{+}(\Sigma p)\left(\xi^{1}, \theta\right) \otimes \Pi_{-}(\Sigma q)$. Since we have changed the orientation, we must change the sign. Therefore:

$$
\begin{aligned}
\int_{\Sigma\left(T^{*} M\right)} & \theta_{1} \wedge \\
= & \theta_{2} \wedge \operatorname{ch}\left(\Pi_{+}(\Sigma r)\right) \\
= & \int_{D_{+} M_{1}} \theta_{1} \wedge \operatorname{ch}\left(\Pi_{+}(\Sigma p)\right) \cdot \int_{\Sigma\left(T^{*} M_{2}\right)} \theta_{2} \wedge \operatorname{ch}\left(\Pi_{+}(\Sigma q)\right) \\
& -\int_{D_{-} M_{1}} \theta_{1} \wedge \operatorname{ch}\left(\Pi_{+}(\Sigma p)\right) \cdot \int_{\Sigma\left(T^{*} M_{2}\right)} \theta_{2} \wedge \operatorname{ch}\left(\Pi_{-}(\Sigma q)\right)
\end{aligned}
$$

$\operatorname{ch}\left(\Pi_{+}(\Sigma q)\right)+c h\left(\Pi_{-}(\Sigma q)\right)=\operatorname{ch}\left(V_{2}\right)$ does not involve the fiber coordinates of $\Sigma\left(T^{*} M_{2}\right)$ and thus

$$
\int_{\Sigma\left(T^{*} M_{2}\right)} \theta_{2} \wedge c h\left(V_{2}\right)=0
$$

We may therefore replace $-\operatorname{ch}\left(\Pi_{-}(\Sigma q)\right)$ by $\operatorname{ch}\left(\Pi_{+}(\Sigma q)\right)$ in evaluating the integral over $D_{-}\left(M_{1}\right) \times \Sigma\left(T^{*} M_{2}\right)$ to complete the proof of (c).

We prove (d) in a similar fashion. We suppose without loss of generality that $p$ and $q$ are homogeneous of degree 1. We parameterize $S\left(T^{*} M\right)=$ $S\left(T^{*} M_{1}\right) \times\left[0, \frac{\pi}{2}\right] \times S\left(T^{*} M_{2}\right)$ in the form $\left(\xi^{1} \cos \theta, \xi^{2} \sin \theta\right)$. Then

$$
r=\left(\begin{array}{cc}
\sin \theta & 0 \\
0 & -\sin \theta
\end{array}\right) \otimes p\left(\xi^{2}\right)+\left(\begin{array}{cc}
0 & (\cos \theta) q^{*}\left(\xi^{1}\right) \\
(\cos \theta) q\left(\xi^{1}\right) & 0
\end{array}\right)
$$

Again we decompose $V_{2}=\Pi_{+}(p) \oplus \Pi_{-}(p)$ so that

$$
\begin{array}{ll}
\Sigma r=\left((\sin \theta) \gamma_{1}+(\cos \theta) \alpha^{1}\left(\xi^{1}\right)\right) & \text { on } V \otimes \Pi_{+}(p) \\
\Sigma r=\left(\sin (-\theta) \gamma_{1}+(\cos \theta) \alpha^{1}\left(\xi^{1}\right)\right) & \text { on } V \otimes \Pi_{-}(p)
\end{array}
$$

where

$$
\gamma_{1}=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right) \quad \alpha^{1}=\left(\begin{array}{rr}
0 & q^{*} \\
q & 0
\end{array}\right) .
$$

The remainder of the argument is exactly as before (with the appropriate change in notation from (c)) and is therefore omitted.

We now check a specific example to verify some normalizing constants:
Lemma 3.9.4. Let $M=S^{1}$ be the unit circle. Define $P: C^{\infty}(M) \rightarrow$ $C^{\infty}(M)$ by:

$$
P\left(e^{i n \theta}\right)= \begin{cases}n e^{i(n-1) \theta} & \text { for } n \geq 0 \\ n e^{i n \theta} & \text { for } n \leq 0\end{cases}
$$

then $P$ is an elliptic pseudo-differential operator with $\operatorname{index}(P)=1$. Furthermore:

$$
\int_{\Sigma\left(T^{*} S^{1}\right)} \operatorname{ch}\left(\Pi_{+}(\Sigma p)\right)=-1 .
$$

Proof: Let $P_{0}=-i \partial / \partial \theta$ and let $P_{1}=\left\{-\partial^{2} / \partial \theta^{2}\right\}^{1 / 2} . P_{0}$ is a differential operator while $P_{1}$ is a pseudo-differential operator by the results of section 1.10. It is immediate that:

$$
\begin{array}{ll}
\sigma_{L}\left(P_{0}\right)=\xi, & P_{0}\left(e^{i n \theta}\right)=n e^{i n \theta} \\
\sigma_{L}\left(P_{1}\right)=|\xi|, & P_{1}\left(e^{i n \theta}\right)=|n| e^{i n \theta}
\end{array}
$$

We define:

$$
\begin{aligned}
Q_{0} & =\frac{1}{2} e^{-i \theta}\left(P_{0}+P_{1}\right) \\
\sigma_{L} Q_{0} & = \begin{cases}\xi e^{-i \theta} & \xi \geq 0 \\
0 & \xi \leq 0\end{cases} \\
Q_{0}\left(e^{i n \theta}\right) & = \begin{cases}n e^{i(n-1) \theta} & n \geq 0 \\
0 & n \leq 0\end{cases}
\end{aligned}
$$

and

$$
\begin{aligned}
Q_{1} & =\frac{1}{2}\left(P_{0}-P_{1}\right) \\
\sigma_{L} Q_{1} & = \begin{cases}0 & \xi \geq 0 \\
\xi & \xi \leq 0\end{cases} \\
Q_{1}\left(e^{i n \theta}\right) & = \begin{cases}0 & n \geq 0 \\
n e^{i n \theta} & n \leq 0\end{cases}
\end{aligned}
$$

Consequently, $P=Q_{0}+Q_{1}$ is a pseudo-differential operator and

$$
\sigma_{L} P=p= \begin{cases}\xi e^{-i \theta} & \xi>0 \\ \xi & \xi<0\end{cases}
$$

It is clear $P$ is surjective so $\mathrm{N}\left(P^{*}\right)=0$. Since $\mathrm{N}(P)$ is the space of constant functions, $\mathrm{N}(P)$ is one dimensional so index $(P)=1$. We compute:

$$
\Sigma p(\theta, \xi, t)= \begin{cases}\left(\begin{array}{cc}
t & \xi e^{+i \theta} \\
\xi e^{-i \theta} & -t
\end{array}\right) & \text { if } \xi \geq 0 \\
\left(\begin{array}{cc}
t & \xi \\
\xi & -t
\end{array}\right) & \text { if } \xi \leq 0\end{cases}
$$

Since $\Sigma p$ does not depend on $\theta$ for $\xi \leq 0$, we may restrict to the region $\xi \geq 0$ in computing the integral. (We must smooth out the symbol to be smooth where $\xi=0$ but suppress these details to avoid undue technical complications).

It is convenient to introduce the parameters:

$$
u=t, \quad v=\xi \cos \theta, \quad w=\xi \sin \theta \quad \text { for } u^{2}+v^{2}+w^{2}=1
$$

then this parametrizes the region of $\Sigma\left(T^{*} S^{1}\right)$ where $\xi \geq 0$ in a 1-1 fashion except where $\xi=0$. Since $d u \wedge d v \wedge d w=-\xi d \theta \wedge d \xi \wedge d t, S^{2}$ inherits the reversed orientation from its natural one. Let

$$
e_{0}=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right), \quad e_{1}=\left(\begin{array}{rr}
0 & 1 \\
1 & 0
\end{array}\right), \quad e_{2}=\left(\begin{array}{rr}
0 & i \\
-i & 0
\end{array}\right)
$$

so $\Sigma p(u, v, w)=u e_{0}+v e_{1}+w e_{2}$. Then by Lemma 2.1.5 we have

$$
-\int_{S^{2}} \operatorname{ch}\left(\Pi_{+}(\Sigma p)\right)=(-1) \cdot\left(\frac{i}{2}\right) \cdot \operatorname{Tr}\left(e_{0} e_{1} e_{2}\right)=-1
$$

which completes the proof.
We can now state the Atiyah-Singer index theorem:

Theorem 3.9.5 (The index theorem). Let $P: C^{\infty}\left(V_{1}\right) \rightarrow C^{\infty}\left(V_{2}\right)$ be an elliptic pseudo-differential operator. Let $T d_{r}(M)=T d(T M \otimes \mathbf{C})$ be the Todd class of the complexification of the real tangent bundle. Then:

$$
\operatorname{index}(P)=(-1)^{\operatorname{dim} M} \int_{\Sigma\left(T^{*} M\right)} T d_{r}(M) \wedge \operatorname{ch}\left(\Pi_{+}(\Sigma p)\right)
$$

We remark that the additional factor of $(-1)^{\operatorname{dim} M}$ could have been avoided if we changed the orientation of $\Sigma\left(T^{*} M\right)$.

We begin the proof by reducing to the case $\operatorname{dim} M=m$ even and $M$ orientable. Suppose first that $m$ is odd. We take $Q: C^{\infty}\left(S^{1}\right) \rightarrow C^{\infty}\left(S^{1}\right)$ to be the operator defined in Lemma 3.9.4 with index +1 . We then form the operator $R$ with index $(R)=\operatorname{index}(P) \operatorname{index}(Q)=\operatorname{index}(P)$ defined in Lemma 3.9.3. Although $R$ is not a pseudo-differential operator, it can be uniformly approximated by pseudo-differential operators in the natural Fredholm topology (once the order of $Q$ is adjusted suitably). (This process does not work when discussing the twisted eta invariant and will involve us in additional technical complications in the next chapter). Therefore:

$$
\begin{aligned}
& (-1)^{m+1} \int_{\Sigma\left(T^{*}\left(M \times S^{1}\right)\right)} T d_{r}\left(M \times S^{1}\right) \wedge \operatorname{ch}\left(\Pi_{+}(\Sigma r)\right) \\
& =(-1)^{m} \int_{\Sigma\left(T^{*} M\right)} T d_{r}(M) \wedge \operatorname{ch}\left(\Pi_{+}(\Sigma p)\right) \cdot(-1) \int_{\Sigma\left(T^{*} S^{1}\right)} \operatorname{ch}\left(\Pi_{+}(\Sigma q)\right) \\
& =(-1)^{m} \int_{\Sigma\left(T^{*} M\right)} T d_{r}(M) \wedge \operatorname{ch}\left(\Pi_{+}(\Sigma p)\right) .
\end{aligned}
$$

To show the last integral gives index $(P)$ it suffices to show the top integral gives index $(R)$ and therefore reduces the proof of Theorem 3.9.3 to the case $\operatorname{dim} M=m$ even.

If $M$ is not orientable, we let $M^{\prime}$ be the orientable double cover of $M$. It is clear the formula on the right hand side of the equation multiplies by two. More careful attention to the methods of the heat equation for pseudo-differential operators gives a local formula for the index even in this case as the left hand side is also multiplied by two under this double cover.

This reduces the proof of Theorem 3.9.3 to the case $\operatorname{dim} M=m$ even and $M$ orientable. We fix an orientation on $M$ henceforth.

Lemma 3.9.6. Let $P: C^{\infty}\left(\Lambda^{+}\right) \rightarrow C^{\infty}\left(\Lambda^{-}\right)$be the operator of the signature complex. Let $\omega=c h_{m / 2}\left(\Pi_{+}(\Sigma p)\right) \in H^{m}\left(\Sigma\left(T^{*} M\right) ; \mathbf{C}\right)$. Then if $\omega_{M}$ is the orientation class of $M$
(a) $\omega_{M} \wedge \omega$ gives the orientation of $\Sigma\left(T^{*} M\right)$.
(b) If $S^{m}$ is a fiber sphere of $\Sigma\left(T^{*} M\right)$, then $\int_{S^{m}} \omega=2^{m / 2}$.

Proof: Let $\left(x_{1}, \ldots, x_{m}\right)$ be an oriented local coordinate system on $M$ so that the $\left\{d x_{j}\right\}$ are orthonormal at $x_{0} \in M$. If $\xi=\left(\xi, \ldots, \xi_{m}\right)$ are the dual fiber coordinates for $T^{*} M$ ), then:

$$
p(\xi)=\sum_{j} i \xi_{j}\left(c\left(d x_{j}\right)\right)=\sum_{j} i \xi_{j}\left(\operatorname{ext}\left(d x_{j}\right)-\operatorname{int}\left(\xi_{j}\right)\right)
$$

gives the symbol of $(d+\delta) ; c(\cdot)$ denotes Clifford multiplication as defined previously. We let $e_{j}=i c\left(d x_{j}\right)$; these are self-adjoint matrices such that $e_{j} e_{k}+e_{k} e_{j}=2 \delta_{j k}$. The orientation class is defined by:

$$
e_{0}=i^{m / 2} c\left(d x_{1}\right) \ldots c\left(d x_{m}\right)=(-i)^{m / 2} e_{1} \ldots e_{m}
$$

The bundles $\Lambda^{ \pm}$are defined as the $\pm 1$ eigenspaces of $e_{0}$. Consequently,

$$
\Sigma p(\xi, t)=t e_{0}+\sum \xi_{j} e_{j}
$$

Therefore by Lemma 2.1.5 when $S^{m}$ is given its natural orientation,

$$
\begin{aligned}
\int_{S^{m}} c h_{m / 2} \Pi_{+}(\Sigma p) & =i^{m / 2} 2^{-m / 2} \operatorname{Tr}\left(e_{0} e_{1} \ldots e_{m}\right) \\
& =i^{m / 2} 2^{-m / 2} \operatorname{Tr}\left(e_{0} i^{m / 2} e_{0}\right) \\
& =(-1)^{m / 2} 2^{-m / 2} \operatorname{Tr}(I)=(-1)^{m / 2} 2^{-m / 2} 2^{m} \\
& =(-1)^{m / 2} 2^{m / 2}
\end{aligned}
$$

However, $S^{m}$ is in fact given the orientation induced from the orientation on $\Sigma\left(T^{*} M\right)$ and on $M$. At the point $(x, 0, \ldots, 0,1)$ in $T^{*} M \oplus \mathbf{R}$ the natural orientations are:

$$
\begin{aligned}
\text { of } X: & d x_{1} \wedge \cdots \wedge d x_{m} \\
\text { of } \Sigma\left(T^{*} M\right): & d x_{1} \wedge d \xi_{1} \wedge \cdots \wedge d x_{m} \wedge d \xi_{m} \\
& =(-1)^{m / 2} d x_{1} \wedge \cdots \wedge d x_{m} \wedge d \xi_{1} \wedge \cdots \wedge d \xi_{m} \\
\text { of } S^{m}: & (-1)^{m / 2} d \xi_{1} \wedge \cdots \wedge d \xi_{m}
\end{aligned}
$$

Thus with the induced orientation, the integral becomes $2^{m / 2}$ and the lemma is proved.

Consequently, $\omega$ provides a cohomology extension and:
Lemma 3.9.7. Let $\rho: \Sigma\left(T^{*} M\right) \rightarrow M$ where $M$ is orientable and even dimensional. Then
(a) $\rho^{*}: H^{*}(M ; \mathbf{C}) \rightarrow H^{*}\left(\Sigma\left(T^{*} M\right) ; \mathbf{C}\right)$ is injective.
(b) If $\omega$ is as defined in Lemma 3.9.6, then we can express any $\alpha \in$ $H^{*}\left(\Sigma\left(T^{*} M\right) ; \mathbf{C}\right)$ uniquely as $\alpha=\rho^{*} \alpha_{1}+\rho^{*} \alpha_{2} \wedge \omega$ for $\alpha_{i} \in H^{*}(M ; \mathbf{C})$.

Since $\rho^{*}$ is injective, we shall drop it henceforth and regard $H^{*}(M ; \mathbf{C})$ as being a subspace of $H^{*}\left(\Sigma\left(T^{*} M\right)\right.$; $\left.\mathbf{C}\right)$.

The Chern character gives an isomorphism $K(X ; \mathbf{C}) \simeq H^{\mathrm{e}}(X ; \mathbf{C})$. When we interpret Lemma 3.9.7 in $K$-theory, we conclude that we can decompose $K\left(\Sigma\left(T^{*} M\right) ; \mathbf{C}\right)=K(M ; \mathbf{C}) \oplus K(M ; \mathbf{C}) \otimes \Pi_{+}(\Sigma p) ; c h(V)$ generates $H^{e}(M ; \mathbf{C})$ as $V$ ranges over $K(M ; \mathbf{C})$. Therefore $K\left(\Sigma\left(T^{*} M\right) ; \mathbf{C}\right) / K(M ; \mathbf{C})$ is generated as an additive module by the twisted signature complex with coefficients in bundles over $M . \Pi_{+}\left(\Sigma p_{V}\right)=V \otimes \Pi_{+}(\Sigma p)$ if $p_{V}$ is the symbol of the signature complex with coefficients in $V$.

In Lemma 3.9.2, we interpreted the index as a map in $K$-theory. Since it is linear, it suffices to compute on the generators given by the signature complex with coefficients in $V$. This proves:

Lemma 3.9.8. Assume $M$ is orientable and of even dimension $m$. Let $P_{V}$ be the operator of the signature complex with coefficients in $V$. The bundles $\left\{\Pi_{+}(\Sigma p)_{V}\right\}_{V \in \operatorname{Vect}(M)}$ generate $K\left(\Sigma\left(T^{*} M\right) ; \mathbf{C}\right) / K(M ; \mathbf{C})$ additively. It suffices to prove Theorem 3.9.5 in general by checking it on the special case of the operators $P_{V}$.

We will integrate out the fiber coordinates to reduce the integral of Theorem 3.9.5 from $\Sigma\left(T^{*} M\right)$ to $M$. We proceed as follows. Let $W$ be an oriented real vector bundle of fiber dimension $k+1$ over $M$ equipped with a Riemannian inner product. Let $S(W)$ be the unit sphere bundle of $W$. Let $\rho: S(W) \rightarrow M$ be the natural projection map. We define a map $\mathcal{I}: C^{\infty}(\Lambda(S(W))) \rightarrow C^{\infty}(\Lambda(M))$ which is a $C^{\infty}(\Lambda(M))$ module homomorphism and which commutes with integration-i.e., if $\alpha \in C^{\infty}(\Lambda(S(W)))$ and $\beta \in C^{\infty}(\Lambda(M))$, we require the map $\alpha \mapsto \mathcal{I}(\alpha)$ to be linear so that $\mathcal{I}\left(\rho^{*} \beta \wedge \alpha\right)=\beta \wedge \mathcal{I}(\alpha)$ and $\int_{S(W)} \alpha=\int_{M} \mathcal{I}(\alpha)$.

We construct $\mathcal{I}$ as follows. Choose a local orthonormal frame for $W$ to define fiber coordinates $u=\left(u_{0}, \ldots, u_{k}\right)$ on $W$. This gives a local representation of $S(W)=V \times S^{k}$ over the coordinate patch $\mathcal{U}$ on $M$. If $\alpha \in C^{\infty}(\Lambda(S(W)))$ has support over $V$, we can decompose $\alpha=\sum_{\nu} \beta_{\nu} \wedge \alpha_{\nu}$ for $\beta_{\nu} \in C^{\infty}(\Lambda(\mathcal{U}))$ and $\alpha_{\nu} \in C^{\infty}\left(\Lambda\left(S^{k}\right)\right)$. We permit the $\alpha_{\nu}$ to have coefficients which depend upon $x \in \mathcal{U}$. This expression is, of course, not unique. Then $\mathcal{I}(\alpha)$ is necessarily defined by:

$$
\mathcal{I}(\alpha)(x)=\sum \beta_{\nu} \int_{S^{k}} \alpha_{\nu}(x) .
$$

It is clear this is independent of the particular way we have decomposed $\alpha$. If we can show $\mathcal{I}$ is independent of the frame chosen, then this will define $\mathcal{I}$ in general using a partition of unity argument.

Let $u_{i}^{\prime}=a_{i j}(x) u_{j}$ be a change of fiber coordinates. Then we compute:

$$
d u_{i}^{\prime}=a_{i j}(x) d u_{j}+d a_{i j}(x) u_{j}
$$

Clearly if $a$ is a constant matrix, we are just reparamatrizing $S^{k}$ so the integral is unchanged. We fix $x_{0}$ and suppose $a\left(x_{0}\right)=I$. Then over $x_{0}$,

$$
d u_{I}=d v_{I}+\sum_{|I|<|J|} c_{I, J} \wedge d v_{J} \quad \text { where } c_{I, J} \in \Lambda^{|I|-|J|}(M)
$$

To integrate and get an answer different from 0 over $S^{k}$, we must have $|I|=k$ so these error terms integrate to zero and $\mathcal{I}$ is invariantly defined.

We specialize to the case $W=T^{*} M \oplus 1$. The orientation of $M$ induces a natural orientation of $T^{*} M \oplus 1$ as a bundle in such a way as to agree with the orientation of $T^{*} M \oplus 1$ as a topological space. Let $\alpha=\operatorname{ch}\left(\Pi_{+}(\Sigma p)\right)$, so $\mathcal{I}(\alpha) \in C^{\infty}(\Lambda(M))$. If we reverse the orientation of $M$, then we interchange the roles of $\Lambda^{+}$and of $\Lambda^{-}$. This has the effect of replacing the parameter $u$ by $-u$ which is equivalent to reversing the orientation of $T^{*} M \oplus 1$ as a topological space. Since both orientations have been reversed, the orientation of the fiber is unchanged so $\mathcal{I}(\alpha)$ is invariantly defined independent of any local orientation of $M$. It is clear from the definition that $\mathcal{I}(\alpha)$ is a polynomial in the jets of the metric and is invariant under changes of the metric by a constant factor. Therefore Theorems 2.5.6 and Lemma 2.5.3 imply $\mathcal{I}(\alpha)$ is a real characteristic form. By Lemma 3.9.6, we can expand $\mathcal{I}(\alpha)=2^{m / 2}+\cdots$.

We solve the equation:

$$
\left\{\mathcal{I}\left(\operatorname{ch}\left(\Pi_{+}(\Sigma p)\right)\right) \wedge \operatorname{Tod} d(m)\right\}_{m-4 s}=2^{(m-4 s) / 2} L_{s}
$$

recursively to define a real characteristic class we shall call $\operatorname{Todd}(m)$ for the moment. It is clear $\operatorname{Tod} d(m)=1+\cdots$. In this equation, $L_{s}$ is the Hirzebruch genus.
Lemma 3.9.9. Let Todd $(m)$ be the real characteristic class defined above. Then if $P$ is any elliptic pseudo-differential operator,

$$
\operatorname{index}(P)=\int_{\Sigma\left(T^{*} M\right)} \operatorname{Todd}(m) \wedge \operatorname{ch}\left(\Pi_{+}(\Sigma p)\right)
$$

Proof: By Lemma 3.9 .8 it suffices to prove this identity if $P$ is the operator of the twisted signature complex. By Theorem 3.1.5,

$$
\begin{aligned}
\operatorname{index}\left(P_{V}^{\text {signature }}\right) & =\int_{M} \sum_{2 t+4 s=m} \operatorname{ch}(V) 2^{t} \wedge L_{s} \\
& =\int_{M} \operatorname{ch}(V) \wedge \mathcal{I}\left(\operatorname{ch}\left(\Pi_{+}(\Sigma p)\right)\right) \wedge \operatorname{Todd}(m) \\
& =\int_{\Sigma\left(T^{*} M\right)} \operatorname{ch}(V) \wedge \operatorname{ch}\left(\Pi_{+}(\Sigma p)\right) \wedge \operatorname{Todd}(m) \\
& =\int_{\Sigma\left(T^{*} M\right)} \operatorname{Todd}(m) \wedge \operatorname{ch}\left(\Pi_{+}\left(\Sigma p_{V}\right)\right)
\end{aligned}
$$

It is clear $\operatorname{Todd}(m)$ is uniquely determined by Lemma 3.9.9. Both the index of an elliptic operator and the formula of Lemma 3.9.9 are multiplicative with respect to products by Lemma 3.9 .3 so $\operatorname{Todd}(m)$ is a multiplicative characteristic form. We may therefore drop the dependence upon the dimension $m$ and simply refer to Todd. We complete the proof of the Atiyah-Singer index theorem by identifying this characteristic form with the real Todd polynomial of $T(M)$.

We work with the Dolbeault complex instead of with the signature complex since the representations involved are simpler. Let $m$ be even, then $(-1)^{m}=1$. Let $M$ be a holomorphic manifold with the natural orientation. We orient the fibers of $T^{*} M$ using the natural orientation which arises from the complex structure on the fibers. If $\xi$ are the fiber coordinates, this gives the orientation:

$$
d x_{1} \wedge d y_{1} \wedge \cdots \wedge d x_{n} \wedge d y_{n} \wedge d \xi_{1} \wedge \cdots \wedge d \xi_{m} \quad \text { where } m=2 n
$$

This gives the total space $T^{*} M$ an orientation which is $(-1)^{n}$ times the simplectic orientation. Let $S^{m}$ denote a fiber sphere of $\Sigma\left(T^{*} M\right)$ with this orientation and let $q$ be the symbol of the Dolbeault complex. Then:

$$
\operatorname{index}(\bar{\partial})=\int_{\Sigma\left(T^{*} M\right)} \operatorname{Todd} \wedge \operatorname{ch}\left(\Pi_{+} \Sigma q\right)=(-1)^{n} \int_{M} \operatorname{Todd} \wedge \mathcal{I}\left(\operatorname{ch}\left(\Pi_{+} \Sigma q\right)\right)
$$

We define the complex characteristic form

$$
\mathcal{S}=(-1)^{m / 2} \mathcal{I}\left(\operatorname{ch}\left(\Pi_{+}(\Sigma q)\right)\right)
$$

then the Riemann-Roch formula implies that:

$$
\mathcal{S} \wedge \operatorname{Todd}(m)=\operatorname{Todd}\left(T_{c} M\right)
$$

It is convenient to extend the definition of the characteristic form $\mathcal{S}$ to arbitrary complex vector bundles $V$. Let $W=\Lambda^{*}(V)$ and let $q(v): V \rightarrow$ $\operatorname{END}(W)$ be defined by Lemma 3.5.2(a) to be the symbol of $\bar{\partial}+\delta^{\prime \prime}$ if $V=T_{c}(M)$. We let ext: $V \rightarrow \operatorname{HOM}(W, W)$ be exterior multiplication. This is complex linear and we let int be the dual of ext; $\operatorname{int}(\lambda v)=\bar{\lambda} \operatorname{int}(v)$ for $\lambda \in \mathbf{C}$. ext is invariantly defined while int requires the choice of a fiber matric. We let $q=\operatorname{ext}(v)-\operatorname{int}(v)$ (where we have deleted the factor of $i / 2$ which appears in Lemma 3.5.2 for the sake of simplicity).

We regard $q$ as a section to the bundle $\operatorname{HOM}(V, \operatorname{HOM}(W, W))$. Fix a Riemannian connection on $V$ and covariantly differentiate $q$ to compute $\nabla q \in C^{\infty}\left(T^{*} M \otimes \operatorname{HOM}(V, \operatorname{HOM}(W, W))\right)$. Since the connection is Riemannian, $\nabla q=0$; this is not true in general for non-Riemannian connections.

If $V$ is trivial with flat connection, the bundles $\Pi_{ \pm}(\Sigma q)$ have curvature $\pi_{ \pm} d \pi_{ \pm} d \pi_{ \pm}$as computed in Lemma 2.1.5. If $V$ is not flat, the curvature of
$V$ enters into this expression. The connection and fiber metric on $V$ define a natural metric on $T^{*} V$. We use the splitting defined by the connection to decompose $T^{*} V$ into horizontal and vertical components. These components are orthogonal with respect to the natural metric on $V$. Over $V, q$ becomes a section to the bundle $\operatorname{HOM}(V, V)$. We let $\nabla^{V}$ denote covariant differentiation over $V$, then $\nabla^{V} q$ has only vertical components in $T^{*} V$ and has no horizontal components.

The calculation performed in Lemma 2.1.5 shows that in this more general setting that the curvatures of $\nabla_{ \pm}$are given by:

$$
\Omega_{ \pm}=\pi_{+}\left(\nabla^{V} \pi_{+} \wedge \nabla^{V} \pi_{+}+\rho^{*} \Omega_{W}\right)
$$

If we choose a frame for $V$ and $W$ which is covariant constant at a point $x_{0}$, then $\nabla^{V} \pi_{+}=d \pi_{+}$has only vertical components while $\rho^{*} \Omega_{W}$ has only horizontal components. If $\Omega_{V}$ is the curvature of $V$, then $\Omega_{W}=\Lambda\left(\Omega_{V}\right)$.

Instead of computing $\mathcal{S}$ on the form level, we work with the corresponding invariant polynomial.

Lemma 3.9.10. Let $A$ be an $n \times n$ complex skew-adjoint matrix. Let $B=\Lambda(A)$ acting on $\Lambda\left(\mathbf{C}^{n}\right)=\mathbf{C}^{2^{n}}$. Define:

$$
\mathcal{S}(A)=\sum_{\nu}(-1)^{n}\left(\frac{i}{2 \pi}\right)^{\nu} \frac{1}{\nu!} \int_{S^{2 n}} \operatorname{Tr}\left\{\left(\pi_{+} d \pi_{+} d \pi_{+}+B\right)^{\nu}\right\}
$$

If $x_{j}=i \lambda_{j} / 2 \pi$ are the normalized eigenvectors of $A$, then

$$
\mathcal{S}(A)=\prod_{j} \frac{e^{x_{j}}-1}{x_{j}}
$$

Proof: If $V=V_{1} \oplus V_{2}$ and if $A=A_{1} \oplus A_{2}$, then the Dolbeault complex decomposes as a tensor product by Lemma 3.9.3. The calculations of Lemma 3.9.3 using the decomposition of the bundles $\Pi_{ \pm}$shows $\mathcal{S}(A)$ is a multiplicative characteristic class. To compute the generating function, it suffices to consider the case $n=1$.

If $n=1, A=\lambda$ so that $B=\left(\begin{array}{ll}0 & 0 \\ 0 & \lambda\end{array}\right)$, if we decompose $W=\Lambda^{0,0} \oplus \Lambda^{0,1}=$ $1 \oplus V$. If $x+i y$ give the usual coordinates on $V=\mathbf{C}$, then:

$$
q(x, y)=x\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right)+y\left(\begin{array}{rr}
0 & -i \\
i & 0
\end{array}\right)
$$

by Lemma 3.5 .2 which gives the symbol of the Dolbeault complex. Therefore:

$$
q(x, y, u)=x\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right)+y\left(\begin{array}{rr}
0 & -i \\
i & 0
\end{array}\right)+u\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right)=x e_{0}+y e_{1}+u e_{2}
$$

We compute:

$$
e_{0} e_{1} e_{2}=i\left(\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right) \quad \text { and } \quad \pi_{+} d \pi_{+} d \pi_{+}=\frac{i}{2} \pi_{+} \mathrm{dvol}
$$

Since $n=1,(-1)^{n}=-1$ and:

$$
\begin{aligned}
\mathcal{S}(A) & =\sum_{j>0}-\frac{i}{2}\left(\frac{i}{2 \pi}\right)^{j} \frac{1}{j!} \int \operatorname{Tr}\left(\left(\pi_{+} \mathrm{dvol}+\pi_{+} B\right)^{j}\right) \\
& =\sum_{j>0}-\frac{i}{2}\left(\frac{i}{2 \pi}\right)^{j} \frac{1}{(j-1)!} \int \operatorname{Tr}\left(\pi_{+} B\right)^{j-1} \mathrm{dvol}
\end{aligned}
$$

We calculate that:

$$
\begin{aligned}
\pi_{+} B & =\frac{1}{2}\left(\begin{array}{cc}
1+u & x-i y \\
x+i y & 1-u
\end{array}\right)\left(\begin{array}{cc}
0 & 0 \\
0 & \lambda
\end{array}\right)=\frac{1}{2}\left(\begin{array}{cc}
0 & * \\
0 & (1-u) \lambda
\end{array}\right) \\
\left(\pi_{+} B\right)^{j-1} & =2^{1-j} \lambda^{j-1}(1-u)^{j-1}\left(\begin{array}{ll}
0 & * \\
0 & 1
\end{array}\right)
\end{aligned}
$$

where "*" indicates a term we are not interested in. We use this identity and re-index the sum to express:

$$
\mathcal{S}(A)=\sum_{j \geq 0} \frac{1}{4 \pi}\left(\frac{i \lambda}{2 \pi}\right)^{j} \frac{1}{j!} 2^{-j} \int_{S^{2}}(1-u)^{j} \mathrm{dvol}
$$

We introduce the integrating factor of $e^{-r^{2}}$ to compute:

$$
\begin{aligned}
\int_{R^{3}} u^{2 k} e^{-r^{2}} d x d y d u & =\pi \int_{R} u^{2 k} e^{-u^{2}} d u \\
& =\int_{0}^{\infty} r^{2 k+2} e^{-r^{2}} d r \cdot \int_{S^{2}} u^{2 k} \mathrm{dvol} \\
& =(2 k+1) / 2 \int_{0}^{\infty} r^{2 k} e^{-r^{2}} d r \cdot \int_{S^{2}} u^{2 k} \mathrm{dvol}
\end{aligned}
$$

so that:

$$
\int_{S^{2}} u^{2 k} \mathrm{dvol}=4 \pi /(2 k+1)
$$

The terms of odd order integrate to zero so:

$$
\begin{aligned}
\int_{S^{2}}(1-u)^{j} \mathrm{dvol} & =4 \pi \sum\binom{j}{2 k} \cdot \frac{1}{2 k+1}=4 \pi \int_{0}^{1} \sum\binom{j}{2 k} t^{2 k} d t \\
& =2 \pi \int_{0}^{1}(1+t)^{j}+(1-t)^{j} d t \\
& =\left.2 \pi \frac{(1+t)^{j+1}-(1-t)^{j+1}}{j+1}\right|_{0} ^{1} \\
& =4 \pi \cdot \frac{2^{j}}{j+1}
\end{aligned}
$$

We substitute this to conclude:

$$
\mathcal{S}(A)=\sum_{j \geq 0}\left(\frac{i \lambda}{2 \pi}\right)^{j} \frac{1}{(j+1)!}
$$

If we introduce $x=i \lambda / 2 \pi$ then

$$
\mathcal{S}(x)=\sum_{j \geq 0} \frac{x^{j}}{(j+1)!}=\frac{e^{x}-1}{x}
$$

which gives the generating function for $\mathcal{S}$. This completes the proof of the lemma.

We can now compute Todd. The generating function of $\operatorname{Todd}\left(T_{c}\right)$ is $x /\left(1-e^{-x}\right)$ so that $\operatorname{Tod} d=\mathcal{S}^{-1} \cdot \operatorname{Todd}\left(T_{c}\right)$ will have generating function:

$$
\frac{x}{1-e^{-x}} \cdot \frac{x}{e^{x}-1}=\frac{x}{1-e^{-x}} \cdot \frac{-x}{1-e^{x}}
$$

which is, of course, the generating function for the real Todd class. This completes the proof. We have gone into some detail to illustrate that it is not particularly difficult to evaluate the integrals which arise in applying the index theorem. If we had dealt with the signature complex instead of the Dolbeault complex, the integrals to be evaluated would have been over $S^{4}$ instead of $S^{2}$ but the computation would have been similar.

## CHAPTER 4

## GENERALIZED

INDEX THEOREMS AND SPECIAL TOPICS

## Introduction

This chapter is less detailed than the previous three as several lengthy calculations are omitted in the interests of brevity. In sections 4.1 through 4.6, we sketch the interrelations between the Atiyah-Patodi-Singer twisted index theorem, the Atiyah-Patodi-Singer index theorem for manifolds with boundary, and the Lefschetz fixed point formulas.

In section 4.1, we discuss the absolute and relative boundary conditions for the de Rham complex if the boundary is non-empty. We discuss Poincaré duality and the Hodge decomposition theorem in this context. The spin, signature, and Dolbeault complexes do not admit such local boundary conditions and the index theorem in this context looks quite different. In section 4.2, we prove the Gauss-Bonnet theorem for manifolds with boundary and identify the invariants arising from the heat equation with these integrands.

In section 4.3, we introduce the eta invariant as a measure of spectral asymmetry and establish the regularity at $s=0$. We discuss without proof the Atiyah-Patodi-Singer index theorem for manifolds with boundary. The eta invariant enters as a new non-local ingredient which was missing in the Gauss-Bonnet theorem. In section 4.4, we review secondary characteristic classes and sketch the proof of the Atiyah-Patodi-Singer twisted index theorem with coefficients in a locally flat bundle. We discuss in some detail explicit examples on a 3 -dimensional lens space.

In section 4.5, we turn to the Lefschetz fixed point formulas. We treat the case of isolated fixed points in complete detail in regard to the four classical elliptic complexes. We also return to the 3 -dimensional examples discussed in section 4.4 to relate the Lefschetz fixed point formulas to the twisted index theorem using results of Donnelly. We discuss in some detail the Lefschetz fixed point formulas for the de Rham complex if the fixed point set is higher dimensional. There are similar results for both the spin and signature complexes which we have omitted for reasons of space. In section 4.6 we use these formulas for the eta invariant to compute the $K$-theory of spherical space forms.

In section 4.7, we turn to a completely new topic. In a lecture at M.I.T., Singer posed the question:

Suppose $P(G)$ is a scalar valued invariant of the metric so that $P(M)=$ $\int_{M} P(G)$ dvol is independent of the metric. Then is there a universal constant $c$ so $P(M)=c \chi(M)$ ?
The answer to this question (and to related questions involving form valued invariants) is yes. This leads immediately to information regarding the higher order terms in the expansion of the heat equation. In section 4.8, we use the functorial properties of the invariants to compute $a_{n}(x, P)$ for an arbitrary second order elliptic parital differential operator with leading symbol given by the metric tensor for $n=0,2,4$. We list (without proof) the corresponding formula if $n=6$. This leads to Patodi's formula for $a_{n}\left(x, \Delta_{p}^{m}\right)$ discussed in Theorem 4.8.18. In section 4.9 we discuss some results of Ikeda to give examples of spherical space forms which are isospectral but not diffeomorphic. We use the eta invariant to show these examples are not even equivariant cobordant.

The historical development of these ideas is quite complicated and in particular the material on Lefschetz fixed point formulas is due to a number of authors. We have included a brief historical survey at the beginning of section 4.6 to discuss this material.

### 4.1. The de Rham Complex for Manifolds with Boundary.

In section 1.9, we derived a formula for the index of a strongly elliptic boundary value problem. The de Rham complex admits suitable boundary conditions leading to the relative and absolute cohomolgy groups. It turns out that the other 3 classical elliptic complexes do not admit even a weaker condition of ellipticity.

Let $(d+\delta): C^{\infty}(\Lambda(M)) \rightarrow C^{\infty}(\Lambda(M))$ be the de Rham complex. In the third chapter we assumed the dimension $m$ of $M$ to be even, but we place no such restriction on the parity here. $M$ is assumed to be compact with smooth boundary $d M$. Near the boundary, we introduce coordinates $(y, r)$ where $y=\left(y_{1}, \ldots, y_{m-1}\right)$ give local coordinates for $d M$ and such that $M=\{x: r(x) \geq 0\}$. We further normalize the coordinates by assuming the curves $x(r)=\left(y_{0}, r\right)$ are unit speed geodesics perpendicular to $d M$ for $r \in[0, \delta)$.

Near $d M$, we decompose any differential form $\theta \in \Lambda(M)$ as

$$
\theta=\theta_{1}+\theta_{2} \wedge d r \quad \text { where } \theta_{i} \in \Lambda(d M)
$$

are tangential differential forms. We use this decomposition to define:

$$
\alpha(\theta)=\theta_{1}-\theta_{2} \wedge d r
$$

$\alpha$ is self-adjoint and $\alpha^{2}=1$. We define the absolute and relative boundary conditions

$$
B_{a}(\theta)=\theta_{2} \quad \text { and } \quad B_{r}(\theta)=\theta_{1}
$$

We let $B$ denote either $B_{a}$ or $B_{r}$, then $B:\left.\Lambda(M)\right|_{d M} \rightarrow \Lambda(d M)$. We note that $B_{a}$ can be identified with orthogonal projection on the -1 eigenspace of $\alpha$ while $B_{r}$ can be identified with orthogonal projection on the +1 eigenspace of $\alpha$. There is a natural inclusion map $i: d M \rightarrow M$ and $B_{r}(\theta)=$ $i^{*}(\theta)$ is just the pull-back of $\theta$. The boundary condition $B_{r}$ does not depend on the Riemannian metric chosen, while the boundary condition $B_{a}$ does depend on the metric.

Lemma 4.1.1. Let $B=B_{a}$ or $B_{r}$, then $(d+\delta, B)$ is self-adjoint and elliptic with respect to the cone $\mathbf{C}-\mathbf{R}_{+}-\mathbf{R}_{-}$.
Proof: We choose a local orthonormal frame $\left\{e_{0}, \ldots, e_{m-1}\right\}$ for $T^{*} M$ near $d M$ so that $e_{0}=d r$. Let

$$
p_{j}=i e_{j}
$$

act on $\Lambda(M)$ by Clifford multiplication. The $p_{j}$ are self-adjoint and satisfy the commutation relation:

$$
p_{j} p_{k}+p_{k} p_{j}=2 \delta_{j k}
$$

The symbol of $(d+\delta)$ is given by:

$$
p(x, \xi)=z p_{0}+\sum_{j=1}^{m-1} \zeta_{j} p_{j}
$$

As in Lemma 1.9.5, we define:

$$
\tau(y, \zeta, \lambda)=i p_{0}\left(\sum_{j=1}^{m-1} \zeta_{j} p_{j}-\lambda\right), \quad(\zeta, \lambda) \neq(0,0) \in T^{*}(d M) \times\left\{\mathbf{C}-\mathbf{R}_{+}-\mathbf{R}_{-}\right\}
$$

We define the new matrices:

$$
q_{0}=p_{0} \quad \text { and } \quad q_{j}=i p_{0} p_{j} \quad \text { for } 1 \leq j \leq m-1
$$

so that:

$$
\tau(y, \zeta, \lambda)=-i \lambda q_{0}+\sum_{j=1}^{m-1} \zeta_{j} q_{j}
$$

The $\left\{q_{j}\right\}$ are self-adjoint and satisfy the commutation relations $q_{j} q_{k}+$ $q_{k} q_{j}=2 \delta_{j k}$. Consequently:

$$
(y, \zeta, \lambda)^{2}=\left(|\zeta|^{2}-\lambda^{2}\right) I
$$

We have $\left(|\zeta|^{2}-\lambda^{2}\right) \in \mathbf{C}-\mathbf{R}_{-}-0$ so we can choose $\mu^{2}=\left(|\zeta|^{2}-\lambda^{2}\right)$ with $\operatorname{Re}(\mu)>0$. Then $\tau$ is diagonalizable with eigenvalues $\pm \mu$ and $V_{ \pm}(\tau)$ is the span of the eigenvectors of $\tau$ corresponding to the eigenvalue $\pm \mu$. (We set $V=\Lambda(M)=\Lambda\left(T^{*} M\right)$ to agree with the notation of section 9.1).

We defined $\alpha\left(\theta_{1}+\theta_{2} \wedge d r\right)=\theta_{1}-\theta_{2} \wedge d r$. Since Clifford multiplication by $e_{j}$ for $1 \leq j \leq m-1$ preserves $\Lambda(d M)$, the corresponding $p_{j}$ commute with $\alpha$. Since Clifford multiplication by $d r$ interchanges the factors of $\Lambda(d M) \oplus \Lambda(d M) \wedge d r, \quad \alpha$ anti-commutes with $p_{0}$. This implies $\alpha$ anticommutes with all the $q_{j}$ and consequently anti-commutes with $\tau$. Thus the only common eigenvectors must belong to the eigenvalue 0 . Since 0 is not an eigenvalue,

$$
V_{ \pm}(\alpha) \cap V_{ \pm}(\tau)=\{0\}
$$

Since $B_{a}$ and $B_{r}$ are just orthogonal projection on the $\mp 1$ eigenspaces of $\alpha, \mathrm{N}\left(B_{a}\right)$ and $\mathrm{N}\left(B_{r}\right)$ are just the $\pm 1$ eigenspaces of $\alpha$. Thus

$$
B: V_{ \pm}(\tau) \rightarrow \Lambda(d M)
$$

is injective. Since $\operatorname{dim}\left(V_{ \pm}(\tau)\right)=\operatorname{dim} \Lambda(d M)=2^{m-1}$, this must be an isomorphism which proves the ellipticity. Since $p_{0}$ anti-commutes with $\alpha$,

$$
p_{0}: V_{ \pm}(\alpha) \rightarrow V_{\mp}(\alpha)
$$

so $(d+\delta)$ is self-adjoint with respect to either $B_{a}$ or $B_{r}$ by Lemma 1.9.5.
By Lemma 1.9.1, there is a spectral resolution for the operator $(d+\delta)_{B}$ in the form $\left\{\lambda_{\nu}, \phi_{\nu}\right\}_{\nu=1}^{\infty}$ where $(d+\delta) \phi_{\nu}=\lambda_{\nu} \phi_{\nu}$ and $B \phi_{\nu}=0$. The $\phi_{\nu} \in C^{\infty}(\Lambda(M))$ and $\left|\lambda_{\nu}\right| \rightarrow \infty$. We set $\Delta=(d+\delta)^{2}=d \delta+\delta d$ and define:

$$
\begin{aligned}
\mathrm{N}\left((d+\delta)_{B}\right) & =\left\{\phi \in C^{\infty}(\Lambda(M)): B \phi=(d+\delta) \phi=0\right\} \\
\mathrm{N}\left((d+\delta)_{B}\right)_{j} & =\left\{\phi \in C^{\infty}\left(\Lambda^{j}(M)\right): B \phi=(d+\delta) \phi=0\right\} \\
\mathrm{N}\left(\Delta_{B}\right) & =\left\{\phi \in C^{\infty}(\Lambda(M)): B \phi=B(d+\delta) \phi=\Delta \phi=0\right\} \\
\mathrm{N}\left(\Delta_{B}\right)_{j} & =\left\{\phi \in C^{\infty}\left(\Lambda^{j}(M)\right): B \phi=B(d+\delta) \phi=\Delta \phi=0\right\}
\end{aligned}
$$

Lemma 4.1.2. Let $B$ denote either the relative or the absolute boundary conditions. Then
(a) $\mathrm{N}\left((d+\delta)_{B}\right)=\mathrm{N}\left(\Delta_{B}\right)$.
(b) $\mathrm{N}\left((d+\delta)_{B}\right)_{j}=\mathrm{N}\left(\Delta_{B}\right)_{j}$.
(c) $\mathrm{N}\left((d+\delta)_{B}\right)=\bigoplus_{j} \mathrm{~N}\left((d+\delta)_{B}\right)_{j}$.

Proof: Let $B \phi=B(d+\delta) \phi=\Delta \phi=0$. Since $(d+\delta)_{B}$ is self-adjoint with respect to the boundary condition $B$, we can compute that $\Delta \phi \cdot \phi=$ $(d+\delta) \phi \cdot(d+\delta) \phi=0$ so $(d+\delta) \phi=0$. This shows $\mathrm{N}\left(\Delta_{B}\right) \subset \mathrm{N}\left((d+\delta)_{B}\right)$ and $\mathrm{N}\left(\Delta_{B}\right)_{j} \subset \mathrm{~N}\left((d+\delta)_{B}\right)_{j}$. The reverse inclusions are immediate which proves (a) and (b). It is also clear that $\mathrm{N}\left((d+\delta)_{B}\right)_{j} \subset \mathrm{~N}\left((d+\delta)_{B}\right)$ for each $j$. Conversely, let $\theta \in \mathrm{N}\left((d+\delta)_{B}\right)$. We decompose $\theta=\theta_{0}+\cdots+\theta_{m}$ into homogeneous pieces. Then $\Delta \theta=\sum_{j} \Delta \theta_{j}=0$ implies $\Delta \theta_{j}=0$ for each $j$. Therefore we must check $B \theta_{j}=B(d+\delta) \theta_{j}=0$ since then $\theta_{j} \in \mathrm{~N}\left(\Delta_{B}\right)_{j}=$ $\mathrm{N}\left((d+\delta)_{B}\right)_{j}$ which will complete the proof.

Since $B$ preserves the grading, $B \theta_{j}=0$. Suppose $B=B_{r}$ is the relative boundary conditions so $B\left(\alpha_{1}+\alpha_{2} \wedge d r\right)=\left.\alpha_{1}\right|_{d M}$. Then $B d=d B$ so $B(d+$ $\delta) \theta=d B(\theta)+B \delta(\theta)=B \delta(\theta)=0$. Since $B$ preserves the homogeneity, this implies $B \delta \theta_{j}=0$ for each $j$. We observed $B d \theta=d B \theta=0$ so $B d \theta_{j}=0$ for each $j$ as well. This completes the proof in this case. If $B=B_{a}$ is the absolute boundary condition, use a similar argument based on the identity $B \delta=\delta B$.

We illustrate this for $m=1$ by considering $M=[0,1]$. We decompose $\theta=f_{0}+f_{1} d x$ to decompose $C^{\infty}(\Lambda(M))=C^{\infty}(M) \oplus C^{\infty}(M) d x$. It is immediate that:

$$
(d+\delta)\left(f_{0}, f_{1}\right)=\left(-f_{1}^{\prime}, f_{0}^{\prime}\right)
$$

so $(d+\delta) \theta=0$ implies $\theta$ is constant. $B_{a}$ corresponds to Dirichlet boundary conditions on $f_{1}$ while $B_{r}$ corresponds to Dirichlet boundary conditions on $f_{0}$. Therefore:

$$
\begin{aligned}
H_{a}^{0}([0,1] ; \mathbf{C}) & =\mathbf{C} & H_{a}^{1}([0,1] ; \mathbf{C})=0 \\
H_{r}^{0}([0,1] ; \mathbf{C}) & =0 & H_{r}^{1}([0,1] ; \mathbf{C})=\mathbf{C}
\end{aligned}
$$

A priori, the dimensions of the vector spaces $H_{a}^{j}(M ; \mathbf{C})$ and $H_{r}^{j}(M ; \mathbf{C})$ depend on the metric. It is possible, however, to get a more invariant definition which shows in fact they are independent of the metric. Lemma 1.9.1 shows these spaces are finite dimensional.

Let $d: C^{\infty}\left(\Lambda^{j}\right) \rightarrow C^{\infty}\left(\Lambda^{j+1}\right)$ be the de Rham complex. The relative boundary conditions are independent of the metric and are $d$-invariant. Let $C_{r}^{\infty}\left(\Lambda^{p}\right)=\left\{\theta \in C^{\infty} \Lambda^{p}: B_{r} \theta=0\right\}$. There is a chain complex $d: C_{r}^{\infty}\left(\Lambda^{p}\right) \rightarrow C_{r}^{\infty}\left(\Lambda^{p+1}\right) \rightarrow \cdots$. We define $H_{r}^{p}(M ; \mathbf{C})=(\text { ker } d / \text { image } d)_{p}$ on $C_{r}^{\infty}\left(\Lambda^{p}\right)$ to be the cohomology of this chain complex. The de Rham theorem for manifolds without boundary generalizes to identify these groups with the relative simplicial cohomology $H^{p}(M, d M ; \mathbf{C})$. If $\theta \in \mathrm{N}\left((d+\delta)_{B_{r}}\right)_{p}$ then $d \theta=B_{r} \theta=0$ so $\theta_{j} \in H^{p}(M, d M ; \mathbf{C})$. The Hodge decomposition theorem discussed in section 1.5 for manifolds without boundary generalizes to identify $H_{r}^{p}(M ; \mathbf{C})=\mathrm{N}\left(\Delta_{B_{r}}\right)_{p}$. If we use absolute boundary conditions and the operator $\delta$, then we can define $H_{a}^{p}(M ; \mathbf{C})=\mathrm{N}\left(\Delta_{B_{a}}\right)_{p}=H^{p}(M ; \mathbf{C})$. We summarize these results as follows:

Lemma 4.1.3. (Hodge decomposition theorem). There are natural isomorphisms between the harmonic cohomology with absolute and relative boundary conditions and the simplicial cohomology groups of $M$ :

$$
H_{a}^{p}(M ; \mathbf{C})=\mathrm{N}\left((d+\delta)_{B_{a}}\right)_{j} \simeq H^{p}(M ; \mathbf{C})
$$

and

$$
H_{r}^{p}(M ; \mathbf{C})=\mathrm{N}\left((d+\delta)_{B_{r}}\right)_{j} \simeq H^{p}(M, d M ; \mathbf{C})
$$

If $M$ is oriented, we let $*$ be the Hodge operator $*: \Lambda^{p} \rightarrow \Lambda^{m-p}$. Since * interchanges the decomposition $\Lambda\left(T^{*} d M\right) \oplus \Lambda\left(T^{*} d M\right) \wedge d r$, it anti-commutes with $\alpha$ and therefore $B_{a}(\theta)=0$ if and only if $B_{r}(* \theta)=0$. Since $d \theta=0$ if and only if $\delta * \theta=0$ and similarly $\delta \theta=0$ if and only if $d * \theta=0$, we conclude:

Lemma 4.1.4. Let $M$ be oriented and let * be the Hodge operator. Then * induces a map, called Poincaré duality,

$$
*: H^{p}(M ; \mathbf{C}) \simeq H_{a}^{p}(M ; \mathbf{C}) \stackrel{\simeq}{\rightarrow} H_{r}^{m-p}(M, \mathbf{C}) \simeq H^{m-p}(M, d M ; \mathbf{C})
$$

We define the Euler-Poincaré characteristics by:

$$
\begin{aligned}
\chi(M) & =\sum(-1)^{p} \operatorname{dim} H^{p}(M ; \mathbf{C}) \\
\chi(d M) & =\sum(-1)^{p} \operatorname{dim} H^{p}(d M ; \mathbf{C}) \\
\chi(M, d M) & =\sum(-1)^{p} \operatorname{dim} H^{p}(M, d M ; \mathbf{C}) .
\end{aligned}
$$

The long exact sequence in cohomology:

$$
\cdots H^{p}(d M ; \mathbf{C}) \leftarrow H^{p}(M ; \mathbf{C}) \leftarrow H^{p}(M, d M ; \mathbf{C}) \leftarrow H^{p-1}(d M ; \mathbf{C}) \cdots
$$

shows that:

$$
\chi(M)=\chi(d M)+\chi(M, d M)
$$

If $m$ is even, then $\chi(d M)=0$ as $d M$ is an odd dimensional manifold without boundary so $\chi(M)=\chi(M, d M)$. If $m$ is odd and if $M$ is orientable, then $\chi(M)=-\chi(M, d M)$ by Poincaré duality. $\chi(M)$ is the index of the de Rham complex with absolute boundary conditions; $\chi(M, d M)$ is the index of the de Rham complex with relative boundary conditions. By Lemma 1.9.3, there is a local formula for the Euler-Poincaré characteristic. Since $\chi(M)=-\chi(M, d M)$ if $m$ is odd and $M$ is orientable, by passing to the double cover if necessary we see $\chi(M)=-\chi(M, d M)$ in general if $m$ is odd. This proves:

Lemma 4.1.5.
(a) If $m$ is even, $\chi(M)=\chi(M, d M)$ and $\chi(d M)=0$.
(b) If $m$ is odd, $\chi(M)=-\chi(M, d M)=\frac{1}{2} \chi(d M)$.

In contrast to the situation of manifolds without boundary, if we pass to the category of manifolds with boundary, there exist non-zero index problems in all dimensions $m$.

In the next subsection, we will discuss the Gauss-Bonnet formula for manifolds with boundary. We conclude this subsection with a brief discussion of the more general ellipticity conditions considered by Atiyah and Bott. Let $Q: C^{\infty}\left(V_{1}\right) \rightarrow C^{\infty}\left(V_{2}\right)$ be an elliptic differential operator of order $d>0$ on the interior-i.e., if $q(x, \xi)$ is the leading symbol of $Q$, then $q(x, \xi): V_{1} \rightarrow V_{2}$ is an isomorphism for $\xi \neq 0$. Let $W_{1}=\left.V_{1} \otimes 1^{d}\right|_{d M}$ be the bundle of Cauchy data. We assume $\operatorname{dim} W_{1}$ is even and let $W_{1}^{\prime}$ be a bundle over $d M$ of dimension $\frac{1}{2}\left(\operatorname{dim} W_{1}\right)$. Let $B: C^{\infty}\left(W_{1}\right) \rightarrow C^{\infty}\left(W_{1}^{\prime}\right)$ be a tangential pseudo-differential operator. We consider the ODE

$$
q\left(y, 0, \zeta, D_{r}\right) f=0, \quad \lim _{r \rightarrow \infty} f(r)=0
$$

and let $V_{+}(\tau)(\zeta)$ be the bundle of Cauchy data of solutions to this equation. We say that $(Q, B)$ is elliptic with respect to the cone $\{0\}$ if for all $\zeta \neq 0$, the map:

$$
\sigma_{g}(B)(y, \zeta): V_{+}(\tau)(\zeta) \rightarrow W_{1}^{\prime}
$$

is an isomorphism (i.e., we can find a unique solution to the ODE such that $\sigma_{g}(B)(y, \zeta) \underline{\gamma} f=f^{\prime}$ is pre-assigned in $\left.W_{1}^{\prime}\right) . V_{+}(\tau)$ is a sub-bundle of $W_{1}$ and is the span of the generalized eigenvectors corresponding to eigenvalues with positive real parts for a suitable endomorphism $\tau(\zeta)$ just as in the first order case. $\sigma_{g}$ is the graded leading symbol as discussed in section 1.9.

This is a much weaker condition than the one we have been considering since the only complex value involved is $\lambda=0$. We study the pair

$$
(Q, B): C^{\infty}\left(V_{1}\right) \rightarrow C^{\infty}\left(V_{2}\right) \oplus C^{\infty}\left(W_{1}^{\prime}\right)
$$

Under the assumption of elliptic with respect to the cone $\{0\}$, this operator is Fredholm in a suitable sense with closed range, finite dimensional nullspace and cokernel. We let index $(Q, B)$ be the index of this problem. The Atiyah-Bott theorem gives a formula for the index of this problem.

There exist elliptic complexes which do not admit boundary conditions satisfying even this weaker notion of ellipticity. Let $q(x, \xi)$ be a first order symbol and expand

$$
q(y, 0, \xi, z)=q_{0} z+\sum_{j=1}^{m-1} q_{j} \zeta_{j}
$$

As in Lemma 1.9.5 we define

$$
\tau=i q_{0}^{-1} \sum_{j=1}^{m-1} q_{j} \zeta_{j}
$$

the ellipticity condition on the interior shows $\tau$ has no purely imaginary eigenvalues for $\zeta \neq 0$. We let $V_{ \pm}(\tau)(\zeta)$ be the sub-bundle of $V$ corresponding to the span of the generalized eigenvectors of $\tau$ corresponding to eigenvalues with positive/negative real part. Then $(Q, B)$ is elliptic if and only if

$$
\sigma_{g}(B)(\zeta): V_{+}(\tau)(\zeta) \rightarrow W_{1}^{\prime}
$$

is an isomorphism for all $\zeta \neq 0$.
Let $S\left(T^{*}(d M)\right)=\left\{\zeta \in T^{*}(d M):|\zeta|^{2}=1\right\}$ be the unit sphere bundle over $d M . \quad V_{ \pm}(\tau)$ define sub-bundles of $V$ over $S\left(T^{*}(d M)\right)$. The existence of an elliptic boundary condition implies these sub-bundles are trivial over the fiber spheres. We study the case in which $q^{*} q=|\zeta|^{2} I$. In this case, $q_{0}^{-1}=q_{0}^{*}$. If we set $p_{j}=i q_{0}^{-1} q_{j}$, then these are self-adjoint and satisfy $p_{j} p_{k}+p_{k} p_{j}=2 \delta_{j k}$. If $m$ is even, then the fiber spheres have dimension $m-2$ which will be even. The bundles $V_{ \pm}(\tau)$ were discussed in Lemma 2.1.5 and in particular are non-trivial if

$$
\operatorname{Tr}\left(p_{1}, \ldots, p_{m-1}\right) \neq 0
$$

For the spin, signature, and Dolbeault complexes, the symbol is given by Clifford multiplication and $p_{1}, \ldots, p_{m-1}$ is multiplication by the orientation form (modulo some normalizing factor of $i$ ). Since the bundles involved were defined by the action of the orientation form being $\pm 1$, this proves:
Lemma 4.1.6. Let $Q: C^{\infty}\left(V_{1}\right) \rightarrow C^{\infty}\left(V_{2}\right)$ denote either the signature, the spin, or the Dolbeault complex. Then there does not exist a boundary condition $B$ so that $(Q, B)$ is elliptic with respect to the cone $\{0\}$.

The difficulty comes, of course, in not permitting the target bundle $W^{\prime}$ to depend upon the variable $\zeta$. In the first order case, there is a natural
pseudo-differential operator $B(\zeta)$ with leading symbol given by projection on $V_{+}(\tau)(\zeta)$. This operator corresponds to global (as opposed to local) boundary conditions and leads to a well posed boundry value problem for the other three classical elliptic complexes. Because the boundary value problem is non-local, there is an additional non-local term which arises in the index theorem for these complexes. This is the eta invariant we will discuss later.

### 4.2. The Gauss-Bonnet Theorem For Manifolds with Boundary.

Let $B$ denote either the absolute or relative boundary conditions for the operator $(d+\delta)$ discussed previously. We let $\chi(M)_{B}$ be either $\chi(M)$ or $\chi(M, d M)$ be the index of the de Rham complex with these boundary conditions. Let $\Delta_{B}^{\text {even }}$ and $\Delta_{B}^{\text {odd }}$ be the Laplacian on even/odd forms with the boundary conditions $B \theta=B(d+\delta) \theta=0$. Let $a_{n}(x, d+\delta)=$ $a_{n}\left(x, \Delta^{\text {even }}\right)-a_{n}\left(x, \Delta^{\text {odd }}\right)$ be the invariants of the heat equation defined in the interior of $M$ which were discussed in Lemma 1.7.4. On the boundary $d M$, let $a_{n}(y, d+\delta, B)=a_{n}\left(y, \Delta_{B}^{\text {even }}\right)-a_{n}\left(y, \Delta_{B}^{\text {odd }}\right)$ be the invariants of the heat equation defined in Lemma 1.9.2. Then Lemma 1.9.3 implies:

$$
\begin{aligned}
\chi(M)_{B}= & \operatorname{Tr}\left\{\exp \left(-t \Delta_{B}^{\text {even }}\right)\right\}-\operatorname{Tr} \exp \left\{\left(-t \Delta_{B}^{\text {odd }}\right)\right\} \\
\sim & \sum_{n=0}^{\infty} t^{(n-m) / 2} \int_{M} a_{n}(x, d+\delta) \operatorname{dvol}(x) \\
& +\sum_{n=0}^{\infty} t^{(n-m+1) / 2} \int_{d M} a_{n}(y, d+\delta, B) \operatorname{dvol}(y) .
\end{aligned}
$$

The interior invariants $a_{n}(x, d+\delta)$ do not depend on the boundary condition so we can apply Lemma 2.4.8 to conclude:

$$
\begin{aligned}
a_{n}(x, d+\delta) & =0 & & \text { if } n<m \text { or if } m \text { is odd } \\
a_{m}(x, d+\delta) & =E_{m} & & \text { is the Euler intergrand if } m \text { is even. }
\end{aligned}
$$

In this subsection we will prove the Gauss-Bonnet theorem for manifolds with boundary and identify the boundary integrands $a_{n}(y, d+\delta, B)$ for $n \leq m-1$.

We let $\mathcal{P}$ be the algebra generated by the $\left\{g_{i j / \alpha}\right\}$ variables for $|\alpha| \neq 0$. We always normalize the coordinate system so $g_{i j}(X, G)\left(x_{0}\right)=\delta_{i j}$. We normalize the coordinate system $x=(y, r)$ near the boundary as discussed in section 4.1; this introduces some additional relations on the $g_{i j / \alpha}$ variables we shall discuss shortly. We let $P(Y, G)\left(y_{0}\right)$ be the evaluation of $P \in \mathcal{P}$ on a metric $G$ and relative to the given coordinate system $Y$ on $d M$. We say that $P$ is invariant if $P(Y, G)\left(y_{0}\right)=P(\bar{Y}, G)\left(y_{0}\right)$ for any two such coordinate systems $Y$ and $\bar{Y}$. We introduce the same notion of homogeneity as that discussed in the second chapter and let $\mathcal{P}_{m, n}^{b}$ be the finite dimensional vector space of invariant polynomials which are homogeneous of order $n$ on a manifold $M$ of dimension $m$. The " $b$ " stands for boundary and emphasizes that these are invariants only defined on $d M$; there is a natural inclusion $\mathcal{P}_{m, n} \rightarrow \mathcal{P}_{m, n}^{b}$; by restricting the admissible coordinate transformations we increase the space of invariants.

Lemma 4.2.1. If $B$ denotes either absolute or relative boundary conditions, then $a_{n}(y, d+\delta, B)$ defines an element of $\mathcal{P}_{m, n}^{b}$.
Proof: By Lemma 1.9.2, $a_{n}(y, d+\delta, B)$ is given by a local formula in the jets of the metric which is invariant. Either by examining the analytic proof of Lemma 1.9.2 in a way similar to that used to prove Lemma 1.7.5 and 2.4.1 or by using dimensional analysis as was done in the proof of Lemma 2.4.4, we can show that $a_{n}$ must be homogeneous of order $n$ and polynomial in the jets of the metric.

Our normalizations impose some additional relations on the $g_{i j / \alpha}$ variables. By hypothesis, the curves $\left(y_{0}, r\right)$ are unit speed geodesics perpendicular to $d M$ at $r=0$. This is equivalent to assuming:

$$
\nabla_{N} N=0 \quad \text { and } \quad g_{j m}(y, 0)=\delta_{j m}
$$

where $N=\partial / \partial r$ is the inward unit normal. The computation of the Christoffel symbols of section 2.3 shows this is equivalent to assuming:

$$
\Gamma_{m m j}=\frac{1}{2}\left(g_{m j / m}+g_{m j / m}-g_{m m / j}\right)=0
$$

If we take $j=m$, this implies $g_{m m / m}=0$. Since $g_{m m}(y, 0) \equiv 1$, we conclude $g_{m m} \equiv 1$ so $g_{m m / \alpha} \equiv 0$. Thus $\Gamma_{m m j}=g_{m j / m}=0$. As $g_{m j}(y, 0)=$ $\delta_{m j}$ we conclude $g_{m j} \equiv \delta_{m j}$ and therefore $g_{j m / \alpha}=0, \quad 1 \leq j \leq m$. We can further normalize the coordinate system $Y$ on $d M$ by assuming $g_{j k / l}(Y, G)\left(y_{0}\right)=0$ for $1 \leq j, k, l \leq m-1$. We eliminate all these variables from the algebra defining $\mathcal{P}$; the remaining variables are algebraically independent.

The only 1-jets of the metric which are left are the $\left\{g_{j k / m}\right\}$ variables for $1 \leq j, k \leq m-1$. The first step in Chapter 2 was to choose a coordinate system in which all the 1-jets of the metric vanish; this proved to be the critical obstruction to studying non-Kaehler holomorphic manifolds. It turns out that the $\left\{g_{j k / m}\right\}$ variables cannot be normalized to zero. They are tensorial and give essentially the components of the second fundamental form or shape operator.

Let $\left\{e_{1}, e_{2}\right\}$ be vector fields on $M$ which are tangent to $d M$ along $d M$. We define the shape operator:

$$
S\left(e_{1}, e_{2}\right)=\left(\nabla_{e_{1}} e_{2}, N\right)
$$

along $d M$. It is clear this expression is tensorial in $e_{1}$. We compute:

$$
\left(\nabla_{e_{1}} e_{2}, N\right)-\left(\nabla_{e_{2}} e_{1}, N\right)=\left(\left[e_{1}, e_{2}\right], N\right)
$$

Since $e_{1}$ and $e_{2}$ are tangent to $d M$ along $d M,\left[e_{1}, e_{2}\right]$ is tangent to $d M$ along $d M$ and thus $\left(\left[e_{1}, e_{2}\right], N\right)=0$ along $d M$. This implies $S\left(e_{1}, e_{2}\right)=$
$S\left(e_{2}, e_{1}\right)$ is tensorial in $e_{2}$. The shape operator defines a bilinear map from $T(d M) \times T(d M) \rightarrow \mathbf{R}$. We compute

$$
\left(\nabla_{\partial / \partial y_{j}} \partial / \partial y_{k}, N\right)=\Gamma_{j k m}=\frac{1}{2}\left(g_{j m / k}+g_{k m / j}-g_{j k / m}\right)=-\frac{1}{2} g_{j k / m}
$$

We can construct a number of invariants as follows: let $\left\{e_{j}\right\}$ be a local orthonormal frame for $T(M)$ such that $e_{m}=N=\partial / \partial r$. Define:

$$
\nabla e_{j}=\sum_{1 \leq k \leq m} \omega_{j k} e_{k} \quad \text { for } \omega_{j k} \in T^{*} M \text { and } \omega_{j k}+\omega_{k j}=0
$$

and

$$
\Omega_{j k}=d \omega_{j k}-\sum_{1 \leq \nu \leq m} \omega_{j \nu} \wedge \omega_{\nu k}
$$

The $\omega_{j m}$ variables are tensorial as $\omega_{j m}=\sum_{k=1}^{m-1} S\left(e_{j}, e_{k}\right) \cdot e^{k}$. We define:

$$
\left.\begin{array}{rl}
Q_{k, m}=c_{k, m} \sum \varepsilon\left(i_{1}, \ldots, i_{m-1}\right) \Omega_{i_{1} i_{2}} & \wedge \cdots \wedge \Omega_{i_{2 k-1}, i_{2 k}} \\
& \wedge \omega_{i_{2 k+1}, m}
\end{array}\right) \cdots \wedge \omega_{i_{m-1}, m} \in \Lambda^{m-1} .
$$

for

$$
c_{k, m}=\frac{(-1)^{k}}{\pi^{p} k!2^{k+p} \cdot 1 \cdot 3 \cdots(2 p-2 k-1)} \quad \text { where } p=\left[\frac{m}{2}\right]
$$

The sum defining $Q_{k}$ is taken over all permutations of $m-1$ indices and defines an $m-1$ form over $M$. If $m$ is even, we define:

$$
E_{m}=\frac{(-1)^{p}}{\left\{2^{m} \pi^{p} p!\right\}} \sum \varepsilon\left(i_{1}, \ldots, i_{m}\right) \Omega_{i_{1} i_{2}} \wedge \cdots \wedge \Omega_{i_{m-1} i_{m}}
$$

as the Euler form discussed in Chapter 2.
$Q_{k, m}$ and $E_{m}$ are the SO-invariant forms on $M . E_{m}$ is defined on all of $M$ while $Q_{k, m}$ is only defined near the boundary. Chern noted that if $m$ is even,

$$
E_{m}=-d\left(\sum_{k} Q_{k, m}\right)
$$

This can also be interpreted in terms of the transgression of Chapter 2. Let $\nabla_{1}$ and $\nabla_{2}$ be two Riemannian connections on $T M$. We defined an $m-1$ form $T E_{m}\left(\nabla_{1}, \nabla_{2}\right)$ so that

$$
d T E_{m}\left(\nabla_{1}, \nabla_{2}\right)=E_{m}\left(\nabla_{1}\right)-E_{m}\left(\nabla_{2}\right)
$$

Near $d M$, we split $T(M)=T(d M) \oplus 1$ as the orthogonal complement of the unit normal. We project the Levi-Civita connection on this decomposition,
and let $\nabla_{2}$ be the projected connection. $\nabla_{2}$ is just the sum of the LeviCivita connection of $T(d M)$ and the trivial connection on 1 and is flat in the normal direction. As $\nabla_{2}$ is a direct sum connection, $E_{m}\left(\nabla_{2}\right)=0$. $\nabla_{1}-\nabla_{2}$ is essentially just the shape operator. $T E_{m}=-\sum Q_{k, m}$ and $d T E_{m}=E_{m}\left(\nabla_{1}\right)=E_{m}$. It is an easy exercise to work out the $Q_{k, m}$ using the methods of section 2 and thereby compute the normalizing constants given by Chern.

The Chern-Gauss-Bonnet theorem for manifolds with boundary in the oriented category becomes:

$$
\chi(M)=\int_{M} E_{m}+\int_{d M} \sum_{k} Q_{k, m} .
$$

In the unoriented category, we regard $E_{m} \operatorname{dvol}(x)$ as a measure on $M$ and $\int Q_{k, m} \operatorname{dvol}(y)$ as a measure on $d M$. If $m$ is odd, of course, $\chi(M)=$ $\frac{1}{2} \chi(d M)=\frac{1}{2} \int_{d M} E_{m-1}$ so there is no difficulty with the Chern-GaussBonnet theorem in this case.

We derive the Chern-Gauss-Bonnet theorem for manifolds with boundary from the theorem for manifolds without boundary. Suppose $m$ is even and that the metric is product near the boundary. Let $\bar{M}$ be the double of $M$ then $\chi(\bar{M})=\int \bar{M} E_{m}=2 \int_{M} E_{m}=2 \chi(M)-\chi(d \bar{M})=2 \chi(M)$ so $\chi(M)=\int_{M} E_{m}$. If the metric is not product near the boundary, let $M^{\prime}=d M \times[-1,0] \cup M$ be the manifold $M$ with a collar sewed on. Let $G_{0}$ be the restriction of the metric on $M$ to the boundary and let $G_{0}^{\prime}$ be the product metric on the collar $d M \times[-1,0]$. Using a partition of unity, extend the original metric on $M$ to a new metric which agrees with $G_{0}^{\prime}$ near $d M \times\{-1\}$ which is the boundary of $M^{\prime}$. Then:

$$
\begin{aligned}
\chi(M)=\chi\left(M^{\prime}\right) & =\int_{M_{1}} E_{m}=\int_{M} E_{m}-\int_{d M \times[-1,0]} d\left(\sum_{k} Q_{k, m}\right) \\
& =\int_{M} E_{m}+\int_{d M} \sum_{k} Q_{k, m}
\end{aligned}
$$

by Stoke's theorem; since the $Q_{k}$ vanish identically near $d M \times\{-1\}$ there is no contribution from this component of the boundary of the collar (we change the sign since the orientation of $d M$ as the boundary of $M$ and as the boundary of $d M \times[-1,0]$ are opposite).

We now study the invariants of the heat equation. We impose no restrictions on the dimension $m$. If $M=S^{1} \times M_{1}$ and if $\theta$ is the usual periodic parameter on $S^{1}$, then there is a natural involution on $\Lambda\left(T^{*} M\right)$ given by interchanging $\psi$ with $d \theta \wedge \psi$ for $\psi \in \Lambda\left(M_{1}\right)$. This involution preserves the boundary conditions and the associated Laplacians, but changes the parity
of the factors involved. This shows $a_{n}(y, d+\delta, B)=0$ for such a product metric. We define:

$$
r: \mathcal{P}_{m, n}^{b} \rightarrow \mathcal{P}_{m-1, n}^{b}
$$

to be the dual of the map $M_{1} \rightarrow S^{1} \times M_{1}$. Then algebraically:

$$
r\left(g_{i j / \alpha}\right)= \begin{cases}0 & \text { if } \operatorname{deg}_{1}\left(g_{i j / \alpha}\right) \neq 0 \\ * & \text { if } \operatorname{deg}_{1}\left(g_{i j / \alpha}\right)=0\end{cases}
$$

where "*" is simply a renumbering to shift all the indices down one. (At this stage, it is inconvenient to have used the last index to indicate the normal direction so that the first index must be used to denote the flat index; denoting the normal direction by the last index is sufficiently cannonical that we have not attempted to adopt a different convention despite the conflict with the notation of Chapter 2). This proves:

Lemma 4.2.2. Let $B$ denote either the relative or absolute boundary conditions. Then $a_{n}(y, d+\delta, B) \in \mathcal{P}_{m, n}^{b}$. Furthermore, $r\left(a_{n}\right)=0$ where $r: \mathcal{P}_{m, n}^{b} \rightarrow \mathcal{P}_{m-1, n}^{b}$ is the restriction map.

We can now begin to identify $a_{n}(y, d+\delta, B)$ using the same techniques of invariance theory applied in the second chapter.
Lemma 4.2.3.. Let $P \in \mathcal{P}_{m, n}^{b}$. Suppose that $r(P)=0$. Then:
(a) $P=0$ if $n<m-1$.
(b) If $n=m-1$, then $P$ is a polynomial in the variables $\left\{g_{i j / m}, g_{i j / k l}\right\}$ for $1 \leq i, j, k, l \leq m-1$. Furthermore, $\operatorname{deg}_{j}(A)=2$ for any monomial $A$ of $P$ and for $1 \leq j \leq m-1$.

Proof: As in the proof of Theorem 2.4.7, we shall count indices. Let $P \neq 0$ and let $A$ be a monomial of $P$. Decompose $A$ in the form:

$$
A=g_{u_{1} v_{1} / \alpha_{1}} \ldots g_{u_{k} v_{k} / \alpha_{k}} g_{i_{1} j_{1} / m} \ldots g_{i_{r} j_{r} / m} \quad \text { for }\left|\alpha_{\nu}\right| \geq 2
$$

(We have chosen our coordinate systems so the only non-zero 1-jets are the $g_{i j / m}$ variables.) Since $r(P)=0, \operatorname{deg}_{1}(A) \neq 0$. Since $P$ is invariant, $\operatorname{deg}_{j}(A)>0$ is even for $1 \leq j \leq m-1$. This yields the inequalities:

$$
2 m-2 \leq \sum_{j \leq m-1} \operatorname{deg}_{j}(A) \quad \text { and } \quad r \leq \operatorname{deg}_{m}(A)
$$

From this it follows that:

$$
2 m-2+r \leq \sum_{j} \operatorname{deg}_{j}(A)=2 r+2 k+\sum_{\nu}\left|\alpha_{\nu}\right|+r=2 r+2 k+n .
$$

Since $\left|\alpha_{\nu}\right| \geq 2$ we conclude

$$
2 k \leq \sum_{\nu}\left|\alpha_{\nu}\right|=n-r
$$

We combine these inequalities to conclude $2 m-2+r \leq 2 n+r$ so $n \geq m-1$. This shows $P=0$ if $n<m-1$ which proves (a). If $n=m-1$, all these inequalities must be equalities. $\left|\alpha_{\nu}\right|=2, \operatorname{deg}_{m}(A)=r$, and $\operatorname{deg}_{j}(A)=2$ for $1 \leq j \leq m-1$. Since the index $m$ only appears in $A$ in the $g_{i j / m}$ variables where $1 \leq i, j \leq m-1$, this completes the proof of (b).

Lemma 2.5 .1 only used invariance under the group $\mathrm{SO}(2)$. Since $P$ is invariant under the action of $\mathrm{SO}(m-1)$, we apply the argument used in the proof of Theorem 2.4.7 to choose a monomial $A$ of $P$ of the form:

$$
A=g_{11 / 22} \ldots g_{2 k-1,2 k-1 / k k} g_{k+1, k+1 / m} \ldots g_{m-1, m-1 / m}
$$

where if $k=0$ the terms of the first kind do not appear and if $k=m-1$, the terms of the second kind do not appear. Since $m-1=2 k+r$, it is clear that $r \equiv m-1 \bmod 2$. We denote such a monomial by $A_{k}$. Since $P \neq 0$ implies $c\left(A_{k}, P\right) \neq 0$ for some $k$, we conclude the dimension of the space of such $P$ is at most the cardinality of $\left\{A_{k}\right\}=\left[\frac{m+1}{2}\right]$. This proves:

LEMMA 4.2.4. Let $r: \mathcal{P}_{m, m-1}^{b} \rightarrow \mathcal{P}_{m-1, m-1}^{b}$ be the restriction map defined earlier. Then $\operatorname{dim} \mathrm{N}(r) \leq\left[\frac{m+1}{2}\right]$.

We can now show:
Lemma 4.2.5. Let $P \in \mathcal{P}_{m, m-1}^{b}$ with $r(P)=0$. Let $i: d M \rightarrow M$ be the inclusion map and $i^{*}: \Lambda^{m-1}\left(T^{*} M\right) \rightarrow \Lambda^{m-1}\left(T^{*}(d M)\right)$ be the dual map. Let $*_{m-1}$ be the Hodge operator on the boundary so $*_{m-1}: \Lambda^{m-1}\left(T^{*}(d M)\right)$ $\rightarrow \Lambda^{0}\left(T^{*}(d M)\right)$. Let $\bar{Q}_{k, m}=*_{m-1}\left(i^{*} Q_{k, m}\right)$. Then the $\left\{\bar{Q}_{k, m}\right\}$ form a basis for $\mathrm{N}(r)$ so we can express $P$ as a linear combination of the $\bar{Q}_{k, m}$.

Proof: It is clear $r\left(\bar{Q}_{k, m}\right)=0$ and that these elements are linearly independent. We have $[(m+1) / 2]$ such elements so by Lemma 4.2 .4 they must be a basis for the kernel of $r$. (If we reverse the orientation we change both the sign of $*$ and $Q$ so $\bar{Q}$ is a scalar invariant.)

We note that if $m$ is odd, then $Q_{m-1, m}=c \cdot \sum \varepsilon\left(i_{1}, \ldots, i_{m-1}\right) \Omega_{i_{1} i_{2}} \wedge \cdots \wedge$ $\Omega_{i_{m-1} i_{m-1}}$ is not the Euler form on the boundary since we are using the Levi-Civita connection on $M$ and not the Levi-Civita connection on $d M$. However, $E_{m-1}$ can be expressed in terms of the $\bar{Q}_{k, m}$ in this situation.

Before proceeding to discuss the heat equation, we need a uniqueness theorem:

Lemma 4.2.6. Let $P=\sum_{k} a_{k} Q_{k, m}$ be a linear combination of the $\left\{Q_{k, m}\right\}$. Suppose that $P \neq 0$. Then there exists a manifold $M$ and a metric $G$ so $\int_{d M} P(G)(y) \neq 0$.
Proof: By assumption not all the $a_{k}=0$. Choose $k$ maximal so $a_{k} \neq 0$. Let $n=m-2 k$ and let $M=S^{2 k} \times D^{m-2 k}$ with the standard metric. (If $m-2 k=1$, we let $M=D^{m}$ and choose a metric which is product near $S^{m-1}$.) We let indices $1 \leq i \leq 2 k$ index a frame for $T^{*}\left(S^{2 k}\right)$ and indices $2 k+1 \leq u \leq m$ index a frame for $T^{*}\left(D^{m-2 k}\right)$. Since the metric is product, $\Omega_{i u}=\omega_{i u}=0$ in this situation. Therefore $Q_{j, m}=0$ if $j<k$. Since $a_{j}=0$ for $j>k$ by assumption, we conclude $\int_{d M} P(G)=a_{k} \int_{d M} Q_{k, m}(G)$, so it suffices to show this integral is non-zero. This is immediate if $n=1$ as $Q_{m-1, m}=E_{m-1}$ and $m-1$ is even.

Let $n \geq 2$. We have $Q_{k, m}(G)=E_{2 k}\left(G_{1}\right) \cdot Q_{0, m-2 k}\left(G_{2}\right)$, since the metric is product. Since

$$
\int_{S^{2 k}} E_{2 k}=\chi\left(S^{2 k}\right)=2 \neq 0
$$

we must only show $\int_{d D^{n}} Q_{0, n}$ is non-zero for all $n>1$. Let $\theta$ be a system of local coordinates on the unit sphere $S^{n-1}$ and let $r$ be the usual radial parameter. If $d s_{e}^{2}$ is the Euclidean metric and $d s_{\theta}^{2}$ is the spherical metric, then

$$
d s_{e}^{2}=r^{2} d s_{\theta}^{2}+d r \cdot d r
$$

From the description of the shape operator given previously we conclude that $S=-d s_{\theta}^{2}$. Let $\left\{e_{1}, \ldots, e_{n-1}\right\}$ be a local oriented orthonormal frame for $T\left(S^{n}\right)$, then $\omega_{i n}=-e^{i}$ and therefore $Q_{0, n}=c \cdot \operatorname{dvol}_{n-1}$ where $c$ is a non-zero constant. This completes the proof.

We combine these results in the following Theorem:
Theorem 4.2.7. (Gauss-Bonnet formula for manifolds with boundary).
(a) Let the dimension $m$ be even and let $B$ denote either the relative or the absolute boundary conditions. Let

$$
\begin{aligned}
Q_{k, m}=c_{k, m} \sum \varepsilon\left(i_{1}, \ldots, i_{m-1}\right) \Omega_{i_{1}, i_{2}} \wedge \cdots & \wedge \Omega_{i_{2 k-1}, i_{2 k}} \\
& \wedge \omega_{i_{2 k+1}, m} \wedge \cdots \wedge \omega_{i_{m-1}, m}
\end{aligned}
$$

for $c_{k, m}=(-1)^{k} /\left(\pi^{m / 2} \cdot k!\cdot 2^{k+m / 2} \cdot 1 \cdot 3 \cdots(m-2 k-1)\right)$. Let $\bar{Q}_{k, m}=$ $*\left(Q_{k, m} \mid d M\right) \in \mathcal{P}_{m, m-1}^{b}$. Then:
(i) $a_{n}(x, d+\delta)=0$ for $n<m$ and $a_{n}(y, d+\delta, B)=0$ for $n<m-1$,
(ii) $a_{m}(x, d+\delta)=E_{m}$ is the Euler integrand,
(iii) $a_{m-1}(y, d+\delta, B)=\sum_{k} \bar{Q}_{k, m}$,
(iv) $\chi(M)=\chi(M, d M)=\int_{M} E_{m} \mathrm{dvol}(x)+\sum_{k} \int_{d M} \bar{Q}_{k, m} \mathrm{dvol}(y)$.
(b) Let the dimension $m$ be odd and let $B_{r}$ be the relative and $B_{a}$ the absolute boundary conditions. Then:
(i) $a_{n}(x, d+\delta)=0$ for all $n$ and $a_{n}\left(y, d+\delta, B_{r}\right)=a_{n}\left(y, d+\delta, B_{a}\right)=0$ for $n<m-1$,
(ii) $a_{m-1}\left(y, d+\delta, B_{a}\right)=\frac{1}{2} E_{m-1}$ and $a_{m-1}\left(y, d+\delta, B_{r}\right)=-\frac{1}{2} E_{m-1}$, (iii) $\chi(M)=-\chi(M, d M)=\frac{1}{2} \int_{d M} E_{m-1} \operatorname{dvol}(y)=\frac{1}{2} \chi(d M)$.

This follows immediately from our previous computations, and Lemmas 4.2.5 and 4.2.6.

The Atayah-Bott theorem gives a generalization of the Atiyah-Singer index theorem for index problems on manifolds with boundary. This theorem includes the Gauss-Bonnet theorem as a special case, but does not include the Atiyah-Patodi-Singer index theorem since the signature, spin, and Dolbeault complexes do not admit local boundary conditions of the form we have been discussing. We will discuss this in more detail in subsection 4.5.

### 4.3. The Regularity at $s=0$ of the Eta Invariant.

In this section, we consider the eta invariant defined in section 1.10. This section will be devoted to proving eta is regular at $s=0$. In the next section we will use this result to discuss the twisted index theorem using coefficients in a locally flat bundle. This invariant appears as a boundary correction term in the index theorem for manifolds with boundary.

We shall assume $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ is a self-adjoint elliptic pseudodifferential operator of order $d>0$. We define

$$
\eta(s, P)=\sum_{\lambda_{i}>0}\left(\lambda_{i}\right)^{-s}-\sum_{\lambda_{i}<0}\left(-\lambda_{i}\right)^{-s} \quad \text { for } \operatorname{Re}(s) \gg 0
$$

and use Theorem 1.10.3 to extend $\eta$ meromorphically to the complex plane with isolated simple poles on the real axis. We define

$$
R(P)=d \cdot \operatorname{Res}_{s=0} \eta(s, P)
$$

We will show $R(P)=0$ so $\eta$ is regular at $s=0$. The first step is to show:
Lemma 4.3.1. Let $P$ and $Q$ be self-adjoint elliptic pseudo-differential operators of order $d>0$.
(a) $P \cdot\left(P^{2}\right)^{v}$ is a self-adjoint elliptic pseudo-differential operator for any $v$ and if $2 v+1>0, R(P)=R\left(P \cdot\left(P^{2}\right)^{v}\right)$.
(b) $R(P \oplus Q)=R(P)+R(Q)$.
(c) There is a local formula $a(x, P)$ in the jets of the symbol of $P$ up to order $d$ so that $R(P)=\int_{M} a(x, P)|\operatorname{dvol}(x)|$.
(d) If $P_{t}$ is a smooth 1-parameter family of such operators, then $R\left(P_{t}\right)$ is independent of the parameter $t$.
(e) If $P$ is positive definite, then $R(P)=0$.
(f) $R(-P)=-R(P)$.

Proof: We have the formal identity: $\eta\left(s, P \cdot\left(P^{2}\right)^{v}\right)=\eta((2 v+1) s, P)$. Since we normalized the residue by multiplying by the order of the operator, (a) holds. The fact that $P \cdot\left(P^{2}\right)^{v}$ is again a pseudo-differential operator follows from the work of Seeley. (b) is an immediate consequence of the definition. (c) and (d) were proved in Lemma 1.10.2 for differential operators. The extension to pseudo-differential operators again follows Seeley's work. If $P$ is positive definite, then the zeta function and the eta function coincide. Since zeta is regular at the origin, (e) follows by Lemma 1.10.1. (f) is immediate from the definition.

We note this lemma continues to be true if we assume the weaker condition $\operatorname{det}(p(x, \xi)-i t) \neq 0$ for $(\xi, t) \neq(0,0) \in T^{*} M \times \mathbf{R}$.

We use Lemma 4.3 .1 to interpret $R(P)$ as a map in $K$-theory. Let $S\left(T^{*} M\right)$ be the unit sphere bundle in $T^{*} M$. Let $V$ be a smooth vector bundle over $M$ equipped with a fiber inner product. Let $p: S\left(T^{*} M\right) \rightarrow$
$\operatorname{END}(V)$ be self-adjoint and elliptic; we assume $p(x, \xi)=p^{*}(x, \xi)$ and $\operatorname{det} p(x, \xi) \neq 0$ for $(x, \xi) \in S\left(T^{*} M\right)$. We fix the order of homogeneity $d>0$ and let $p_{d}(x, \xi)$ be the extension of $p$ to $T^{*} M$ which is homogeneous of degree $d$. (In general, we must smooth out the extension near $\xi=0$ to obtain a $C^{\infty}$ extension, but we suppress such details in the interests of notational clarity.)

Lemma 4.3.2. Let $p: S\left(T^{*} M\right) \rightarrow \operatorname{END}(V)$ be self-adjoint and elliptic. Let $d>0$ and let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ have leading symbol $p_{d}$. Then $R(P)$ depends only on $p$ and not on the order $d$ nor the particular operator $P$.

Proof: Let $P^{\prime}$ have order $d$ with the same leading symbol $p_{d}$. We form the elliptic family $P_{t}=t P^{\prime}+(1-t) P$. By Lemma 4.3.1, $R\left(P_{t}\right)$ is independent of $t$ so $R\left(P^{\prime}\right)=R(P)$. Given two different orders, let $(1+2 v) d=d^{\prime}$. Let $Q=P\left(P^{2}\right)^{v}$ then $R(Q)=R(P)$ by Lemma 4.3.1(a). The leading symbol of $Q$ is $p\left(p^{2}\right)^{v} . \quad p^{2}$ is positive definite and elliptic. We construct the homotopy of symbols $q_{t}(x, \xi)=p(x, \xi)\left(t p^{2}(x, \xi)+(1-t)|\xi|^{2}\right)$. This shows the symbol of $Q$ restricted to $S\left(T^{*} M\right)$ is homotopic to the symbol of $P$ restricted to $S\left(T^{*} M\right)$ where the homotopy remains within the class of self-adjoint elliptic symbols. Lemma 4.3 .1 completes the proof.

We let $r(p)=R(P)$ for such an operator $P$. Lemma 4.3.1(d) shows $r(p)$ is a homotopy invariant of $p$. Let $\Pi_{ \pm}(p)$ be the subspaces of $V$ spanned by the eigenvectors of $p(x, \xi)$ corresponding to positive/negative eigenvalues. These bundles have constant rank and define smooth vector bundles over $S\left(T^{*} M\right)$ so $\Pi_{+} \oplus \Pi_{-}=V$. In section 3.9, an essential step in proving the Atayah-Singer index theorem was to interpret the index as a map in $K$-theory. To show $R(P)=0$, we must first interpret it as a map in $K$ theory. The natural space in which to work is $K\left(S\left(T^{*} M\right) ; \mathbf{Q}\right)$ and not $K\left(\Sigma\left(T^{*} M\right) ; \mathbf{Q}\right)$.

Lemma 4.3.3. Let $G$ be an abelian group and let $R(P) \in G$ be defined for any self-adjoint elliptic pseudodifferential operator. Assume $R$ satisfies properties (a), (b), (d), (e) and (f) [but not necessarily (c)] of Lemma 4.3.1. Then there exists a Z-linear map $r: K\left(S\left(T^{*} M\right)\right) \rightarrow G$ so that:
(a) $R(P)=r\left(\Pi_{+}(p)\right)$,
(b) If $\tau: S\left(T^{*} M\right) \rightarrow M$ is the natural projection, then $r\left(\tau^{*} V\right)=0$ for all $V \in K(M)$ so that

$$
r: K\left(S\left(T^{*} M\right)\right) / K(M) \rightarrow G
$$

Remark: We shall apply this lemma in several contexts later, so state it in somewhat greater generality than is needed here. If $G=\mathbf{R}$, we can extend $r$ to a $\mathbf{Q}$ linear map

$$
r: K\left(S\left(T^{*} M\right) ; \mathbf{Q}\right) \rightarrow \mathbf{R} .
$$

Proof: Let $p: S\left(T^{*} M\right) \rightarrow \operatorname{END}(V)$ be self-adjoint and elliptic. We define the bundles $\Pi_{ \pm}(p)$ and let $\pi_{ \pm}(x, \xi)$ denote orthogonal projection on $\Pi_{ \pm}(p)(x, \xi)$. We let $p_{0}=\pi_{+}-\pi_{-}$and define $p_{t}=t p+(1-t) p_{0}$ as a homotopy joining $p$ and $p_{0}$. It is clear that the $p_{t}$ are self-adjoint. Fix $(x, \xi)$ and let $\left\{\lambda_{i}, v_{i}\right\}$ be a spectral resolution of the matrix $p(x, \xi)$. Then:

$$
\begin{aligned}
p(x, \xi)\left(\sum c_{i} v_{i}\right) & =\sum \lambda_{i} c_{i} v_{i} \\
p_{0}(x, \xi)\left(\sum c_{i} v_{i}\right) & =\sum \operatorname{sign}\left(\lambda_{i}\right) c_{i} v_{i}
\end{aligned}
$$

Consequently

$$
p_{t}\left(\sum c_{i} v_{i}\right)=\sum\left(t \lambda_{i}+(1-t) \operatorname{sign}\left(\lambda_{i}\right)\right) c_{i} v_{i}
$$

Since $t \lambda_{i}+(1-t) \operatorname{sign}\left(\lambda_{i}\right) \neq 0$ for $t \in[0,1]$, the family $p_{t}$ is elliptic. We therefore assume henceforth that $p(x, \xi)^{2}=I$ on $S\left(T^{*} M\right)$ and $\pi_{ \pm}=$ $\frac{1}{2}(1 \pm p)$.

We let $k$ be large and choose $V \in \operatorname{Vect}_{k}\left(S\left(T^{*} M\right)\right)$. Choose $W \in$ $\operatorname{Vect}_{k}\left(S\left(T^{*} M\right)\right)$ so $V \oplus W \simeq 1^{2 k}$. We may choose the metric on $1^{2 k}$ so this direct sum is orthogonal and we define $\pi_{ \pm}$to be orthogonal projection on $V$ and $W$. We let $p=\pi_{+}-\pi_{-}$so $\Pi_{+}=V$ and $\Pi_{-}=W$. We define $r(V)=r(p)$. We must show this is well defined on Vect ${ }_{k}$. $W$ is unique up to isomorphism but the isomorphism $V \oplus W=1^{2 k}$ is non-canonical. Let $u: V \rightarrow \bar{V}$ and $v: W \rightarrow \bar{W}$ be isomorphisms where we regard $\bar{V}$ and $\bar{W}$ as orthogonal complements of $1^{2 k}$ (perhaps with another trivialization). We must show $r(p)=r(\bar{p})$. For $t \in[0,1]$ we let

$$
\begin{aligned}
V(t) & =\operatorname{span}(t \cdot v \oplus(1-t) \cdot u(v))_{v \in V} \\
& \subseteq V \oplus W \oplus \bar{V} \oplus \bar{W}=1^{2 k} \oplus 1^{2 k}=1^{4 k}
\end{aligned}
$$

This is a smooth 1-parameter family of bundles connecting $V \oplus 0$ to $0 \oplus \bar{V}$ in $1^{4 k}$. This gives a smooth 1-parameter family of symbols $p(t)$ connecting $p \oplus\left(-1_{2 k}\right)$ to $\left(-1_{2 k}\right) \oplus \bar{p}$. Thus $r(p)=r\left(p \oplus-1_{2 k}\right)=r(p(t))=r\left(-1_{2 k} \oplus \bar{p}\right)=$ $r(\bar{p})$ so this is in fact a well defined map $r: \operatorname{Vect}_{k}\left(S\left(T^{*} M\right)\right) \rightarrow G$. If $V$ is the trivial bundle, then $W$ is the trivial bundle so $p$ decomposes as the direct sum of two self-adjoint matrices. The first is positive definite and the second negative definite so $r(p)=0$ by Lemma 4.3.1(f). It is clear that $r\left(V_{1} \oplus V_{2}\right)=r\left(V_{1}\right)+r\left(V_{2}\right)$ by Lemma 4.3.1(b) and consequently since $r(1)=0$ we conclude $r$ extends to an additive map from $\widetilde{K}\left(S\left(T^{*} M\right)\right) \rightarrow G$. We extend $r$ to be zero on trivial bundles and thus $r: K\left(S\left(T^{*} M\right)\right) \rightarrow G$.

Suppose $V=\tau^{*} V_{0}$ for $V_{0} \in \operatorname{Vect}_{k}(M)$. We choose $W_{0} \in \operatorname{Vect}_{k}(M)$ so $V_{0} \oplus W_{0} \simeq 1^{2 k}$. Then $p=p_{+} \oplus p_{-}$for $p_{+}: S\left(T^{*} M\right) \rightarrow \operatorname{END}\left(V_{0}, V_{0}\right)$
and $p_{-}: S\left(T^{*} M\right) \rightarrow \operatorname{END}\left(W_{0}, W_{0}\right)$. By Lemma 4.3.1, we conclude $r(p)=$ $r\left(p_{+}\right)+r\left(p_{-}\right)$. Since $p_{+}$is positive definite, $r\left(p_{+}\right)=0$. Since $p_{-}$is negative definite, $r\left(p_{-}\right)=0$. Thus $r(p)=0$ and $r\left(\tau^{*} V_{0}\right)=0$.

This establishes the existence of the map $r: K\left(S\left(T^{*} M\right)\right) \rightarrow G$ with the desired properties. We must now check $r(p)=r\left(\Pi_{+}\right)$for general $p$. This follows by definition if $V$ is a trivial bundle. For more general $V$, we first choose $W$ so $V \oplus W=1$ over $M$. We let $q=p \oplus 1$ on $V \oplus W$, then $r(q)=r(p)+r(1)=r(p)$. However, $q$ acts on a trivial bundle so $r(q)=r\left(\Pi_{+}(q)\right)=r\left(\Pi_{+}(p) \oplus \tau^{*} W\right)=r\left(\Pi_{+}(p)\right)+r\left(\tau^{*} W\right)=r\left(\Pi_{+}(p)\right)$ which completes the proof.

Of course, the bundles $\Pi_{ \pm}(p)$ just measure the infinitesimal spectral asymmetry of $P$ so it is not surprising that the bundles they represent in $K$-theory are related to the eta invariant. This construction is completely analogous to the construction given in section 3.9 which interpreted the index as a map in $K$-theory. We will return to this construction again in discussing the twisted index with coefficients in a locally flat bundle.

Such operators arise naturally from considering boundary value problems. Let $M$ be a compact oriented Riemannian manifold of dimension $m=2 k-1$ and let $N=M \times[0,1)$; we let $n \in[0,1)$ denote the normal parameter. Let

$$
(d+\delta)_{+}: C^{\infty}\left(\Lambda^{+}\left(T^{*} N\right)\right) \rightarrow C^{\infty}\left(\Lambda^{-}\left(T^{*} N\right)\right)
$$

be the operator of the signature complex. The leading symbol is given by Clifford multiplication. We can use $c(d n)$, where $c$ is Clifford multiplication, to identify these two bundles over $N$. We express:

$$
(d+\delta)_{+}=c(d n)(\partial / \partial n+A)
$$

The operator $A$ is a tangential differential operator on $C^{\infty}\left(\Lambda^{+}\left(T^{*} N\right)\right)$; it is called the tangential operator of the signature complex. Since we have $c(d n) * c(d n)=-1$, the symbol of $A$ is $-i c(d n) c(\xi)$ for $\xi \in T^{*}(M)$. It is immediate that the leading symbol is self-adjoint and elliptic; since $A$ is natural this implies $A$ is a self-adjoint elliptic partial differential operator on $M$.

Let $\left\{e_{1}, \ldots, e_{m}\right\}$ be a local oriented orthonormal frame for $T^{*} M$. Define:

$$
\omega_{m}=i^{k} e_{1} * \cdots * e_{m} \quad \text { and } \quad \omega_{m+1}=i^{k}(-d n) * e_{1} * \cdots * e_{m}=-d n * \alpha_{m}
$$

as the local orientations of $M$ and $N . \omega_{m}$ is a central element of $\operatorname{CLIF}(M)$; $\omega_{m}^{2}=\omega_{m+1}^{2}=1$. If $\phi \in \Lambda(M)$ define $\tau_{ \pm}(\phi)=\phi \pm c\left(\omega_{m+1}\right) \phi$. This gives an
isomorphism $\tau_{ \pm}: \Lambda(M) \rightarrow \Lambda^{ \pm}(N)$. We compute:

$$
\begin{aligned}
\left(\tau_{+}\right)^{-1}\{-c(d n) c(\xi)\} & \} \tau_{+} \phi \\
& =\left(\tau_{+}\right)^{-1}\left\{-c(d n) c(\xi) \phi-c(d n) c(\xi) c(-d n) c\left(\omega_{m}\right) \phi\right\} \\
& =\left(\tau_{+}\right)^{-1}\left\{-c(d n) c\left(\omega_{m}\right) c\left(\omega_{m}\right) c(\xi)+c\left(\omega_{m}\right) c(\xi) \phi\right\} \\
& =\left(\tau_{+}\right)^{-1}\left\{\left(c\left(\omega_{m+1}\right)+1\right) c\left(\omega_{m}\right) c(\xi) \phi\right\} \\
& =c\left(\omega_{m}\right) c(\xi) \phi
\end{aligned}
$$

If we use $\tau_{+}$to regard $A$ as an operator on $C^{\infty}(\Lambda(M))$ then this shows that $A$ is given by the diagram:

$$
C^{\infty}(\Lambda(M)) \xrightarrow{\nabla} C^{\infty}\left(T^{*}(M \otimes \Lambda(M)) \xrightarrow{c} C^{\infty}(\Lambda(M)) \xrightarrow{\omega} C^{\infty}(\Lambda(M))\right.
$$

where $\omega=\omega_{m}$. This commutes with the operator $(d+\delta) ; \quad A=\omega(d+\delta)$.
Both $\omega$ and $(d+\delta)$ reverse the parity so we could decompose $A=A^{\text {even }} \oplus$ $A^{\text {odd }}$ acting on smooth forms of even and odd degree. We let $p=\sigma_{L}(A)=$ $i c(\omega * \xi)=\sum_{j} \xi_{j} f_{j}$. The $\left\{f_{j}\right\}$ are self-adjoint matrices satisfying the commutation relation

$$
f_{j} f_{k}+f_{k} f_{j}=2 \delta_{j k}
$$

We calculate:

$$
f_{1} \cdots f_{m}=i^{m} c\left(\omega^{m} * e_{1} * \cdots * e_{m}\right)=i^{m-k} c\left(\omega^{m} * \omega\right)=i^{m-k} c(1)
$$

so that if we integrate over a fiber sphere with the natural (not simplectic) orientation,

$$
\begin{aligned}
\int_{S^{m-1}} \operatorname{ch}\left(\Pi_{+}(p)\right) & =i^{k-1} 2^{1-k} \operatorname{Tr}\left(f_{1} \cdots f_{m}\right) \\
& =i^{k-1} 2^{1-k} i^{m-k} 2^{m}=i^{k-1} 2^{1-k} i^{k-1} 2^{2 k-1} \\
& =(-1)^{k-1} \cdot 2^{k}
\end{aligned}
$$

In particular, this is non-zero, so this cohomology class provides a cohomology extension to the fiber.

As a $H^{*}(M ; \mathbf{Q})$ module, we can decompose $H^{*}(S(M) ; \mathbf{Q})=H^{*}(M ; \mathbf{Q}) \oplus$ $x H^{*}(M ; \mathbf{Q})$ where $x=\operatorname{ch}\left(\Pi_{+}(p)\right)$. If we twist the operator $A$ by taking coefficients in an auxilary bundle $V$, then we generate $x H^{*}(M ; \mathbf{Q})$. The same argument as that given in the proof of Lemma 3.9.8 permits us to interpret this in $K$-theory:

Lemma 4.3.4. Let $M$ be a compact oriented Riemannian manifold of dimension $m=2 k-1$. Let $A$ be the tangential operator of the signature complex on $M \times[0,1)$. If $\left\{e_{j}\right\}$ is an oriented local orthonormal basis for $T(M)$, let $\omega=i^{k} e_{1} * \cdots * e_{m}$ be the orientation form acting by Clifford multiplication on the exterior algebra. $A=\omega(d+\delta)$ on $C^{\infty}(\Lambda(M))$. If the symbol of $A$ is $p$,

$$
\int_{S^{m-1}} \operatorname{ch}\left(\Pi_{+}(p)\right)=2^{k}(-1)^{k-1}
$$

The natural map $K(M ; \mathbf{Q}) \rightarrow K\left(S\left(T^{*}(M) ; \mathbf{Q}\right)\right.$ is injective and the group $K\left(S\left(T^{*} M\right) ; \mathbf{Q}\right) / K(M ; \mathbf{Q})$ is generated by the bundles $\left\{\Pi_{+}\left(A_{V}\right)\right\}$ as $V$ runs over $K(M)$.
Remark: This operator can also be represented in terms of the Hodge operator. On $C^{\infty}\left(\Lambda^{2 p}\right)$ for example it is given by $i^{k}(-1)^{p+1}(* d-d *)$. We are using the entire tangential operator (and not just the part acting on even or odd forms). This will produce certain factors of 2 differing from the formulas of Atiyah-Patodi-Singer. If $M$ admits a $\mathrm{SPIN}_{c}$ structure, one can replace $A$ by the tangential operator of the $\mathrm{SPIN}_{c}$ complex; in this case the corresponding integrand just becomes $(-1)^{k-1}$.

We can use this representation to prove:
Lemma 4.3.5. Let $\operatorname{dim} M$ be odd and let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be a selfadjoint elliptic pseudo-differential operator of order $d>0$. Then $R(P)=$ 0 -i.e., $\eta(s, P)$ is regular at $s=0$.

Proof: We first suppose $M$ is orientable. By Lemma 4.3.4, it suffices to prove $R\left(A_{V}\right)=0$ since $r$ is defined in $K$-theory and would then vanish on the generators. However, by Lemma 4.3.1(c), the residue is given by a local formula. The same analysis as that done for the heat equation shows this formula must be homogeneous of order $m$ in the jets of the metric and of the connection on $V$. Therefore, it must be expressible in terms of Pontrjagin forms of $T M$ and Chern forms of $V$ by Theorem 2.6.1. As $m$ is odd, this local formula vanishes and $R\left(A_{V}\right)=0$. If $M$ is not orientable, we pass to the oriented double cover. If $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ over $M$, we let $P^{\prime}: C^{\infty}\left(V^{\prime}\right) \rightarrow C^{\infty}\left(V^{\prime}\right)$ be the lift to the oriented double cover. Then $R(P)=\frac{1}{2} R\left(P^{\prime}\right)$. But $R\left(P^{\prime}\right)=0$ since the double cover is oriented and thus $R(P)=0$. This completes the proof.

This result is due to Atayah, Patodi, and Singer. The trick used in section 3.9 to change the parity of the dimension by taking products with a problem over the circle does not go through without change as we shall see. Before considering the even dimensional case, we must first prove a product formula.

Lemma 4.3.6. Let $M_{1}$ and $M_{2}$ be smooth manifolds. Let $P: C^{\infty}\left(V_{1}\right) \rightarrow$ $C^{\infty}\left(V_{2}\right)$ be an elliptic complex over $M_{1}$. Let $Q: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be a self-adjoint elliptic operator over $M_{2}$. We assume $P$ and $Q$ are differential operators of the same order and form:

$$
R=\left(\begin{array}{cc}
Q & P^{*} \\
P & -Q
\end{array}\right) \quad \text { on } C^{\infty}\left(V_{1} \otimes V \oplus V_{2} \otimes V\right)
$$

then $\eta(s, R)=\operatorname{index}(P) \cdot \eta(s, Q)$.
Proof: This lemma gives the relationship between the index and the eta invariant which we will exploit henceforth. We perform a formal computation. Let $\left\{\lambda_{\nu}, \phi_{\nu}\right\}_{\nu=1}^{\infty}$ be a spectral resolution of the operator $Q$ on $C^{\infty}(V)$. We let $\Delta=P^{*} P$ and decompose $C^{\infty}\left(V_{1}\right)=\mathrm{N}(\Delta)+R(\Delta)$. We let $\left\{\mu_{j}, \theta_{j}\right\}$ be a spectral resolution of $\Delta$ restricted to $\mathrm{N}(\Delta)^{\perp}=\mathrm{R}(\Delta)$. The $\mu_{j}$ are positive real numbers; $\left\{\mu_{j}, P \theta_{j} / \sqrt{\mu_{j}}\right\}$ form a spectral resolution of $\Delta^{\prime}=P P^{*}$ on $\mathrm{N}\left(\Delta^{\prime}\right)^{\perp}=\mathrm{R}\left(\Delta^{\prime}\right)=\mathrm{R}(P)$.

We decompose $L^{2}\left(V_{1} \otimes V\right)=\mathrm{N}(\Delta) \otimes L^{2}(V) \oplus \mathrm{R}(\Delta) \otimes L^{2}(V)$ and $L^{2}\left(V_{2} \otimes\right.$ $V)=\mathrm{N}\left(\Delta^{\prime}\right) \otimes L^{2}(V) \oplus \mathrm{R}\left(\Delta^{\prime}\right) \otimes L^{2}(V) . \operatorname{In} \mathrm{R}(\Delta) \otimes L^{2}(V) \oplus \mathrm{R}\left(\Delta^{\prime}\right) \otimes L^{2}(V)$ we study the two-dimensional subspace that is spanned by the elements: $\left\{\theta_{j} \otimes \phi_{\nu}, P \theta_{j} / \sqrt{\mu_{j}} \otimes \phi_{\nu}\right\}$. The direct sum of these subspaces as $j, \nu$ vary is $\mathrm{R}(\Delta) \otimes L^{2}(V) \oplus \mathrm{R}\left(\Delta^{\prime}\right) \otimes L^{2}(V)$. Each subspace is invariant under the operator $R$. If we decompose $R$ relative to this basis, it is represented by the $2 \times 2$ matrix:

$$
\left(\begin{array}{cc}
\lambda_{\nu} & \sqrt{\mu_{j}} \\
\sqrt{\mu_{j}} & -\lambda_{\nu}
\end{array}\right)
$$

This matrix has two eigenvalues with opposite signs: $\pm \sqrt{\lambda_{\nu}^{2}+\mu_{j}}$. Since $\lambda_{\nu}>0$ these eigenvalues are distinct and cancel in the sum defining eta. Therefore the only contribution to eta comes from $\mathrm{N}(\Delta) \otimes L^{2}(V)$ and $\mathrm{N}\left(\Delta^{\prime}\right) \otimes L^{2}(V)$. On the first subspace, $R$ is $1 \otimes Q$. Each eigenvalue of $Q$ is repeated $\operatorname{dim} \mathrm{N}(\Delta)$ times so the contribution to eta is $\operatorname{dim} \mathrm{N}(\Delta) \eta(s, Q)$. On the second subspace, $R$ is $1 \otimes-Q$ and the contribution to eta is $-\operatorname{dim} \mathrm{N}(\Delta) \eta(s, Q)$. When we sum all these contributions, we conclude that:

$$
\eta(s, R)=\operatorname{dim} \mathrm{N}(\Delta) \eta(s, Q)-\operatorname{dim} \mathrm{N}\left(\Delta^{\prime}\right) \eta(s, Q)=\operatorname{index}(P) \eta(s, Q)
$$

Although this formal cancellation makes sense even if $P$ and $Q$ are not differential operators, $R$ will not be a pseudo-differential operator if $P$ and $Q$ are pseudo-differential operators in general. This did not matter when we studied the index since the index was constant under approximations. The eta invariant is a more delicate invariant, however, so we cannot use the same trick. Since the index of any differential operator on the circle is
zero, we cannot use Lemmas 4.3.5 and 4.3.6 directly to conclude $R(P)=0$ if $m$ is even.

We recall the construction of the operator on $C^{\infty}\left(S^{1}\right)$ having index 1. We fix $a \in \mathbf{R}$ as a real constant and fix a positive order $d \in \mathbf{Z}$. We define:

$$
\begin{aligned}
Q_{0} & =-i \partial / \partial \theta, \quad Q_{1}(a)=\left(Q_{0}^{2 d}+a^{2}\right)^{1 / 2 d} \\
Q_{2}(a) & =\frac{1}{2}\left(e^{-i \theta}\left(Q_{0}+Q_{1}\right)+Q_{0}-Q_{1}\right) \\
Q(a) & =Q_{2}(a) \cdot\left(Q_{1}(a)\right)^{d-1}
\end{aligned}
$$

The same argument as that given in the proof of Lemma 3.9.4 shows these are pseudo-differential operators on $C^{\infty}\left(S^{1}\right) . Q_{2}(a)$ and $Q(a)$ are elliptic families. If $a=0$, then $Q_{2}(0)$ agrees with the operator of Lemma 3.9.4 so index $Q_{2}(0)=1$. Since the index is continuous under perturbation, index $Q_{2}(a)=1$ for all values of $a . \quad Q_{1}(a)$ is self-adjoint so its index is zero. Consequently index $Q(a)=\operatorname{index} Q_{2}(a)+\operatorname{index} Q_{1}(a)=1$ for all $a$.

We let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be an elliptic self-adjoint partial differential operator of order $d>0$ over a manifold $M$. On $M \times S^{1}$ we define the operators:

$$
\begin{aligned}
& Q_{0}=-i \partial / \partial \theta, \quad Q_{1}=\left(Q_{0}^{2 d}+P^{2}\right)^{1 / 2 d} \\
& Q_{2}=\frac{1}{2}\left(e^{-i \theta}\left(Q_{0}+Q_{1}\right)+Q_{0}-Q_{1}\right), \quad Q=Q_{2} \cdot Q_{1}^{d-1}
\end{aligned}
$$

on $C^{\infty}(V)$. These are pseudo-differential operators over $M \times S^{1}$ since $Q_{0}^{2 d}+P^{2}$ has positive definite leading order symbol. We define $R$ by:

$$
R=\left(\begin{array}{cc}
P & Q^{*} \\
Q & -P
\end{array}\right): C^{\infty}(V \oplus V) \rightarrow C^{\infty}(V \oplus V)
$$

This is a pseudo-differential operator of order $d$ which is self-adjoint. We compute:

$$
R^{2}=\left(\begin{array}{cc}
P^{2}+Q^{*} Q & 0 \\
0 & P^{2}+Q Q^{*}
\end{array}\right)
$$

so

$$
\sigma_{L}\left(R^{2}\right)(\xi, z)=\left(\begin{array}{cc}
p(\xi)^{2}+q^{*} q(\xi, z) & 0 \\
0 & p(\xi)^{2}+q q^{*}(\xi, z)
\end{array}\right)
$$

for $\xi \in T^{*} M$ and $z \in T^{*} S^{1}$. Suppose $R$ is not elliptic so $\sigma_{L}(R)(\xi, z) v=0$ for some vector $v \in V \oplus V$. Then $\left\{\sigma_{L}(R)(\xi, z)\right\}^{2} v=0$. We decompose $v=v_{1} \oplus v_{2}$. We conclude:

$$
\left(p(\xi)^{2}+q^{*} q(\xi, z)\right) v_{1}=\left(p(\xi)^{2}+q q^{*}(\xi, z)\right) v_{2}=0
$$

Using the condition of self-adjointness this implies:

$$
p(\xi) v_{1}=q(\xi, z) v_{1}=p(\xi) v_{2}=q^{*}(\xi, z) v_{2}=0
$$

We suppose $v \neq 0$ so not both $v_{1}$ and $v_{2}$ are zero. Since $P$ is elliptic, this implies $\xi=0$. However, for $\xi=0$, the operators $Q_{1}$ and $Q_{2}$ agree with the operators of Lemma 3.9.4 and are elliptic. Therefore $q(\xi, z) v_{1}=$ $q^{*}(\xi, z) v_{2}=0$ implies $z=0$. Therefore the operator $R$ is elliptic.

Lemma 4.3.6 generalizes in this situation to become:
Lemma 4.3.7. Let $P, Q, R$ be defined as above, then $\eta(s, R)=\eta(s, P)$.
Proof: We let $\left\{\lambda_{\nu}, \phi_{\nu}\right\}$ be a spectral resolution of $P$ on $C^{\infty}(V)$ over $M$. This gives an orthogonal direct sum decomposition:

$$
L^{2}(V) \text { over } M \times S^{1}=\bigoplus_{\nu} L^{2}\left(S^{1}\right) \otimes \phi_{\nu}
$$

Each of these spaces is invariant under both $P$ and $R$. If $R_{\nu}$ denotes the restriction of $R$ to this subspace, then

$$
\eta(s, R)=\sum_{\nu} \eta\left(s, R_{\nu}\right)
$$

On $L^{2}\left(\phi_{\nu}\right), \quad P$ is just multiplication by the real eigenvalue $\lambda_{\nu}$. If we replace $P$ by $\lambda_{\nu}$ we replace $Q$ by $Q\left(\lambda_{\nu}\right)$, so $R_{\nu}$ becomes:

$$
R_{\nu}=\left(\begin{array}{cc}
\lambda_{\nu} & Q^{*}\left(\lambda_{\nu}\right) \\
Q\left(\lambda_{\nu}\right) & -\lambda_{\nu}
\end{array}\right) .
$$

We now apply the argument given to prove Lemma 4.3.6 to conclude:

$$
\eta\left(s, R_{\nu}\right)=\operatorname{sign}\left(\lambda_{\nu}\right)\left|\lambda_{\nu}\right|^{-s} \operatorname{index}\left(Q\left(\lambda_{\nu}\right)\right)
$$

Since index $Q\left(\lambda_{\nu}\right)=1$, this shows $\eta(s, R)=\operatorname{sign}\left(\lambda_{\nu}\right)\left|\lambda_{\nu}\right|^{-s}$ and completes the proof.

We can now generalize Lemma 4.3.5 to all dimensions:
Theorem 4.3.8. Let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be a self-adjoint elliptic pseudo-differential operator of order $d>0$. Then $R(P)=0$-i.e., $\eta(s, P)$ is regular at $s=0$.

Proof: This result follows from Lemma 4.3.5 if $\operatorname{dim} M=m$ is odd. If $\operatorname{dim} M=m$ is even and if $P$ is a differential operator, then we form the pseudo-differential operator $P$ over $M \times S^{1}$ with $\eta(s, R)=\eta(s, P)$. Then $\eta(s, R)$ is regular at $s=0$ implies $\eta(s, P)$ is regular at $s=0$. This proves

Theorem 4.3.8 for differential operators. Of course, if $P$ is only pseudodifferential, then $R$ need not be pseudo-differential so this construction does not work. We complete the proof of Theorem 4.3 .8 by showing the partial differential operators of even order generate $K\left(S\left(T^{*} M\right) ; \mathbf{Q}\right) / K(M ; \mathbf{Q})$ if $m$ is even. (We already know the operators of odd order generate if $m$ is odd and if $M$ is oriented).

Consider the involution $\xi \mapsto-\xi$ of the tangent space. This gives a natural $\mathbf{Z}_{2}$ action on $S\left(T^{*} M\right)$. Let $\pi: S\left(T^{*} M\right) \rightarrow S\left(T^{*} M\right) / \mathbf{Z}_{2}=\mathbf{R} P\left(T^{*} M\right)$ be the natural projection on the quotient projective bundle. Since $m$ is even, $m-1$ is odd and $\pi^{*}$ defines an isomorphism between the cohomology of two fibers

$$
\pi^{*}: H^{*}\left(\mathbf{R} P^{m-1} ; \mathbf{Q}\right)=\mathbf{Q} \oplus \mathbf{Q} \rightarrow H^{*}\left(S^{m-1} ; \mathbf{Q}\right)=\mathbf{Q} \oplus \mathbf{Q}
$$

The Kunneth formula and an appropriate Meyer-Vietoris sequence imply that

$$
\pi^{*}: H^{*}\left(\mathbf{R} P\left(T^{*} M\right) ; \mathbf{Q}\right) \rightarrow H^{*}\left(S\left(T^{*} M\right) ; \mathbf{Q}\right)
$$

is an isomorphism in cohomology for the total spaces. We now use the Chern isomorphism between cohomology and $K$-theory to conclude there is an isomorphism in $K$-theory

$$
\pi^{*}: K\left(\mathbf{R} P\left(T^{*} M\right) ; \mathbf{Q}\right) \simeq K\left(S\left(T^{*} M\right) ; \mathbf{Q}\right)
$$

Let $\mathrm{S}_{0}(k)$ be the set of all $k \times k$ self-adjoint matrices $A$ such that $A^{2}=1$ and $\operatorname{Tr}(A)=0$. We noted in section 3.8 that if $k$ is large, $\widetilde{K}(X)=\left[X, \mathrm{~S}_{0}(2 k)\right]$. Thus $\widetilde{K}\left(S\left(T^{*} M\right) ; \mathbf{Q}\right)=\widetilde{K}\left(\mathbf{R} P\left(T^{*} M\right) ; \mathbf{Q}\right)$ is generated by maps $p: S\left(T^{*} M\right) \rightarrow \mathrm{S}_{0}(2 k)$ such that $p(x, \xi)=p(x,-\xi)$. We can approximate any even map by an even polynomial using the StoneWeierstrass theorem. Thus we may suppose $p$ has the form:

$$
p(x, \xi)=\left.\sum_{\substack{|\alpha| \leq n \\|\alpha| \text { even }}} p_{\alpha}(x) \xi^{\alpha}\right|_{S\left(T^{*} M\right)}
$$

where the $p_{\alpha}: M \rightarrow \mathrm{~S}_{0}(2 k)$ and where $n$ is large. As $\alpha$ is even, we can replace $\xi^{\alpha}$ by $\xi^{\alpha}|\xi|^{\{n-|\alpha|\} / 2}$ and still have a polynomial with the same values on $S\left(T^{*} M\right)$. We may therefore assume that $p$ is a homogeneous even polynomial; this is the symbol of a partial differential operator which completes the proof.

If $m$ is odd and if $M$ is orientable, we constructed specific examples of operators generating $K\left(S\left(T^{*} M\right) ; \mathbf{Q}\right) / K(M ; \mathbf{Q})$ using the tangential operator of the signature complex with coefficients in an arbitrary coefficient bundle. If $m$ is even, it is possible to construct explicit second order oprators generating this $K$-theory group. One can then prove directly that
eta is regular at $s=0$ for these operators as they are all "natural" in a certain suitable sense. This approach gives more information by explcitly exhibiting the generators; as the consturction is quite long and technical we have chosen to give an alternate argument based on $K$-theory and refer to (Gilkey, The residue of the global eta function at the origin) for details.

We have given a global proof of Theorem 4.3.8. In fact such a treatment is necessary since in general the local formulas giving the residue at $s=0$ are non-zero. Let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be self-adjoint and elliptic. If $\left\{\lambda_{\nu}, \phi_{\nu}\right\}$ is a spectral resolution of $P$, we define:

$$
\eta(s, P, x)=\sum_{\nu} \operatorname{sign}\left(\lambda_{\nu}\right)\left|\lambda_{\nu}\right|^{-s}\left(\phi_{\nu}, \phi_{\nu}\right)(x)
$$

so that:

$$
\eta(s, P)=\int_{M} \eta(s, P, x) \operatorname{dvol}(x)
$$

Thus $\operatorname{Res}_{s=0} \eta(s, P, x)=a(P, x)$ is given by a local formula.
We present the following example (Gilkey, The residue of the local eta function at the origin) to show this local formula need not vanish identically in general.

Example 4.3.9. Let

$$
e_{1}=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right), \quad e_{2}=\left(\begin{array}{rr}
0 & 1 \\
1 & 0
\end{array}\right), \quad e_{3}=\left(\begin{array}{rr}
0 & i \\
-i & 0
\end{array}\right)
$$

be Clifford matrices acting on $\mathbf{C}^{2}$. Let $T^{m}$ be the $m$-dimensional torus with periodic parameters $0 \leq x_{j} \leq 2$ for $1 \leq j \leq m$.
(a) Let $m=3$ and let $b(x)$ be a real scalar. Let $P=i \sum_{j} e_{j} \partial / \partial x_{j}+b(x) I$. Then $P$ is self-adjoint and elliptic and $a(x, P)=c \Delta b$ where $c \neq 0$ is some universal constant.
(b) Let $m=2$ and let $b_{1}$ and $b_{2}$ be imaginary scalar functions. Let $P=$ $e_{1} \partial^{2} / \partial x_{1}^{2}+e_{2} \partial^{2} / \partial x_{2}^{2}+2 e_{3} \partial^{2} / \partial x_{1} \partial x_{2}+b_{1} e_{1}+b_{2} e_{2}$. Then $P$ is self-adjoint and elliptic and $a(x, P)=c^{\prime}\left(\partial b_{1} / \partial x_{2}-\partial b_{2} / \partial x_{1}\right)$ where $c^{\prime} \neq 0$ is some universal constant.
(c) By twisting this example with a non-trivial index problem and using Lemma 4.3.6, we can construct examples on $T_{m}$ so that $a(x, P)$ does not vanish identically for any dimension $m \geq 2$.

The value of eta at the origin plays a central role in the Atiyah-PatodiSinger index theorem for manifolds with boundary. In section 2, we discussed the transgression briefly. Let $\nabla_{i}$ be two connections on $T(N)$. We defined $T L_{k}\left(\nabla_{1}, \nabla_{2}\right)$ so that $d T L_{k}\left(\nabla_{1}, \nabla_{2}\right)=L_{k}\left(\nabla_{1}\right)-L_{k}\left(\nabla_{2}\right)$; this is a secondary characteristic class. Let $\nabla_{1}$ be the Levi-Civita connection of $N$ and near $M=d N$ let $\nabla_{2}$ be the product connection arising from the
product metric. As $L_{k}\left(\nabla_{2}\right)=0$ we see $d T L_{k}\left(\nabla_{1}, \nabla_{2}\right)=0$. This is the analogous term which appeared in the Gauss-Bonnet theorem; it can be computed in terms of the first and second fundamental forms. For example, if $\operatorname{dim} N=4$

$$
L_{1}(R)=\frac{-1}{24 \cdot \pi^{2}} \operatorname{Tr}(R \wedge R) \quad T L_{1}(R, \omega)=\frac{-1}{24 \cdot 8 \cdot \pi^{2}} \operatorname{Tr}(R \wedge \omega)
$$

where $R$ is the curvature 2 -form and $\omega$ is the second fundamental form.
Theorem 4.3.10 (Atiyah-Patodi-Singer index theorem for MANIFOLDS WITH BOUNDARY). Let $N$ be a $4 k$ dimensional oriented compact Riemannian manifold with boundary M. Then:

$$
\operatorname{signature}(N)=\int_{N} L_{k}-\int_{M} T L_{k}-\frac{1}{2} \eta(0, A)
$$

where $A$ is the tangential operator of the signature complex discussed in Lemma 4.3.4.

In fact, the eta invariant more generally is the boundary correction term in the index theorem for manifolds with boundary. Let $N$ be a compact Riemannian manifold with boundary $M$ and let $P: C^{\infty}\left(V_{1}\right) \rightarrow C^{\infty}\left(V_{2}\right)$ be an elliptic first order differential complex over $N$. We take a metric which is product near the boundary and identify a neighborhood of $M$ in $N$ with $M \times[0,1)$. We suppose $P$ decomposes in the form $P=\sigma(d n)(\partial / \partial n+A)$ on this collared neighborhood where $A$ is a self-adjoint elliptic first order operator over $M$. This is in fact the case for the signature, Dolbeault, or spin compexes. Let $B$ be the spectral projection on the non-negative eigenvalues of $A$. Then:
Theorem 4.3.11. (The Atiyah-Patodi-Singer Index Theorem)
Adopt the notation above. $P$ with boundary condition $B$ is an elliptic problem and

$$
\operatorname{index}(P, B)=\int_{N}\left\{a_{n}\left(x, P^{*} P\right)-a_{n}\left(x, P P^{*}\right)\right\}-\frac{1}{2}\{\eta(0, A)+\operatorname{dim} \mathrm{N}(A)\}
$$

In this expression, $n=\operatorname{dim} N$ and the invariants $a_{n}$ are the invariants of the heat equation discussed previously.

Remark: If $M$ is empty then this is nothing but the formula for index $(P)$ discussed previously. If the symbol does not decompose in this product structure near the boundary of $N$, there are corresponding local boundary correction terms similar to the ones discussed previously. If one takes $P$ to be the operator of the signature complex, then Theorem 4.3.10 can be derived from this more general result by suitably interpreting index $(P, B)=$ signature $(\mathrm{N})-\frac{1}{2} \operatorname{dim} \mathrm{~N}(A)$.

### 4.4. The Eta Invariant with Coefficients <br> In a Locally Flat Bundle.

The eta invariant plays a crucial role in the index theorem for manifolds with boundary. It is also possible to study the eta invariant with coefficients in a locally flat bundle to get a generalization of the Atiyah-Singer theorem.

Let $\rho: \pi_{1}(M) \rightarrow \mathrm{U}(k)$ be a unitary representation of the fundamental group. Let $\bar{M}$ be the universal cover of $M$ and let $\bar{m} \rightarrow g \bar{m}$ for $\bar{m} \in \bar{M}$ and $g \in \pi_{1}(M)$ be the acion of the deck group. We define the bundle $V_{\rho}$ over $M$ by the identification:

$$
V_{\rho}=\bar{M} \times \mathbf{C}^{k} \bmod (\bar{m}, z)=(g \bar{m}, \rho(g) z) .
$$

The transition functions of $V_{\rho}$ are locally constant and $V_{\rho}$ inherits a natural unitary structure and connection $\nabla_{\rho}$ with zero curvature. The holonomy of $\nabla_{\rho}$ is just the representation $\rho$. Conversely, given a unitary bundle with locally constant transition functions, we can construct the connection $\nabla$ to be unitary with zero curvature and recover $\rho$ as the holonomy of the connection. We assume $\rho$ is unitary to work with self-adjoint operators, but all the constructions can be generalized to arbitrary representations in $\mathrm{GL}(k, \mathbf{C})$.

Let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$. Since the transition functions of $V_{\rho}$ are locally constant, we can define $P_{\rho}$ on $C^{\infty}\left(V \otimes V_{\rho}\right)$ uniquely using a partition of unity if $P$ is a differential operator. If $P$ is only pseudo-differential, $P_{\rho}$ is well defined modulo infinitely smoothing terms. We define:

$$
\begin{aligned}
\tilde{\eta}(P) & =\frac{1}{2}\{\eta(P)+\operatorname{dim} \mathrm{N}(P)\} \quad \bmod \mathbf{Z} \\
\operatorname{ind}(\rho, P) & =\tilde{\eta}\left(P_{\rho}\right)-k \tilde{\eta}(P) \quad \bmod \mathbf{Z} .
\end{aligned}
$$

Lemma 4.4.1. $\operatorname{ind}(\rho, P)$ is a homotopy invariant of $P$. If we fix $\rho$, we can interpret this as a map

$$
\operatorname{ind}(\rho, *): K\left(S\left(T^{*} M\right)\right) / K(M) \rightarrow \mathbf{R} \bmod \mathbf{Z}
$$

such that $\operatorname{ind}(\rho, P)=\operatorname{ind}\left(\rho, \Pi_{+}\left(\sigma_{L}(P)\right)\right)$.
Proof: We noted previously that $\tilde{\eta}$ was well defined in $\mathbf{R} \bmod \mathbf{Z}$. Let $P(t)$ be a smooth 1-parameter family of such operators and let $P^{\prime}(t)=\frac{d}{d t} P(t)$. In Theorem 1.10.2, we proved:

$$
\frac{d}{d t} \tilde{\eta}\left(P_{t}\right)=\int_{M} a\left(x, P(t), P^{\prime}(t)\right) \operatorname{dvol}(x)
$$

was given by a local formula. Let $P_{k}=P \otimes 1^{k}$ acting on $V \otimes 1^{k}$. This corresponds to the trivial representation of $\pi_{1}(M)$ in $\mathrm{U}(k)$. The operators
$P_{\rho}$ and $P_{k}$ are locally isomorphic modulo $\infty$ smoothing terms which don't affect the local invariant. Thus $a\left(x, P_{k}(t), P_{k}^{\prime}(t)\right)-a\left(x, P_{\rho}(t), P_{\rho}^{\prime}(t)\right)=0$. This implies $\frac{d}{d t} \operatorname{ind}\left(\rho, P_{t}\right)=0$ and completes the proof of homotopy invariance. If the leading symbol of $P$ is definite, then the value $\eta(0, P)$ is given by a local formula so $\operatorname{ind}(\rho, *)=0$ in this case, as the two local formulas cancel. This verifies properties (d) and (e) of Lemma 4.3.1; properties (a), (b) and (f) are immediate. We therefore apply Lemma 4.3.3 to regard

$$
\operatorname{ind}(\rho, *): K\left(S\left(T^{*} M\right)\right) / K(M) \rightarrow \mathbf{R} \bmod \mathbf{Z}
$$

which completes the proof.
In sections 4.5 and 4.6 we will adopt a slightly different notation for this invariant. Let $G$ be a group and let $R(G)$ be the group representation ring generated by the unitary representations of $G$. Let $R_{0}(G)$ be the ideal of representations of virtual dimension 0 . We extend $\tilde{\eta}: R(G) \rightarrow \mathbf{R} \bmod \mathbf{Z}$ to be a $\mathbf{Z}$-linear map. We let $\operatorname{ind}(\rho, P)$ denote the restriction to $R_{0}(G)$. If $\rho$ is a representation of dimension $j$, then $\operatorname{ind}(\rho, P)=\operatorname{ind}(\rho-j \cdot 1, P)=$ $\eta\left(P_{\rho}\right)-j \eta(P)$. It is convenient to use both notations and the context determines whether we are thinking of virtual representations of dimension 0 or the projection of an actual representation to $R_{0}(G)$. $\operatorname{ind}(\rho, P)$ is not topological in $\rho$, as the following example shows. We will discuss $K$-theory invariants arising from the $\rho$ dependence in Lemma 4.6.5.
Example 4.4.2: Let $M=S^{1}$ be the circle with periodic parameter $0 \leq \theta \leq$ $2 \pi$. Let $g(\theta)=e^{i \theta}$ be the generator of $\pi_{1}(M) \simeq \mathbf{Z}$. We let $\varepsilon$ belong to $\mathbf{R}$ and define:

$$
\rho_{\varepsilon}(g)=e^{2 \pi i \varepsilon}
$$

as a unitary representation of $\pi_{1}(M)$. The locally flat bundle $V_{\rho}$ is topologically trivial since any complex bundle over $S^{1}$ is trivial. If we define a locally flat section to $S^{1} \times \mathbf{C}$ by:

$$
\vec{s}(\theta)=e^{i \varepsilon \theta}
$$

then the holonomy defined by $\vec{s}$ gives the representation $\rho_{\varepsilon}$ since

$$
\vec{s}(2 \pi)=e^{2 \pi i \varepsilon} \vec{s}(0)
$$

We let $P=-i \partial / \partial \theta$ on $C^{\infty}\left(S^{1}\right)$. Then

$$
P_{\varepsilon}=P_{\rho_{\varepsilon}}=e^{+i \varepsilon \theta} P e^{-i \varepsilon \theta}=P-\varepsilon .
$$

The spectrum of $P$ is $\{n\}_{n \in \mathbf{Z}}$ so the spectrum of $P_{\varepsilon}$ is $\{n-\varepsilon\}_{n \in \mathbf{Z}}$. Therefore:

$$
\eta\left(s, P_{\varepsilon}\right)=\sum_{n-\varepsilon \neq 0} \operatorname{sign}(n-\varepsilon)|n-\varepsilon|^{-s} .
$$

We differentiate with respect to $\varepsilon$ to get:

$$
\frac{d}{d \varepsilon}\left(s, P_{\varepsilon}\right)=s \sum_{n-\varepsilon \neq 0}|n-\varepsilon|^{-s-1}
$$

We evaluate at $s=0$. Since the sum defining this shifted zeta function ranges over all integers, the pole at $s=0$ has residue 2 so we conclude:

$$
\frac{d}{d \varepsilon} \tilde{\eta}\left(P_{\varepsilon}\right)=2 \cdot \frac{1}{2}=1 \quad \text { and } \quad \operatorname{ind}\left(\rho_{\varepsilon}, P\right)=\int_{0}^{\varepsilon} 1 d \varepsilon=\varepsilon
$$

We note that if we replace $\varepsilon$ by $\varepsilon+j$ for $j \in \mathbf{Z}$, then the representation is unchanged and the spectrum of the operator $P_{\varepsilon}$ is unchanged. Thus reduction $\bmod \mathbf{Z}$ is essential in making $\operatorname{ind}\left(\rho_{\varepsilon}, P\right)$ well defined in this context.

If $V_{\rho}$ is locally flat, then the curvature of $\nabla_{\rho}$ is zero so $\operatorname{ch}\left(V_{\rho}\right)=0$. This implies $V_{\rho}$ is a torsion element in $K$-theory so $V_{\rho} \otimes 1^{n} \simeq 1^{k n}$ for some integer $n$. We illustrate this with
Example 4.4.3: Let $M=\mathbf{R} P_{3}=S^{3} / \mathbf{Z}_{2}=\mathrm{SO}(3)$ be real projective space in dimension 3 so $\pi_{1}(M)=\mathbf{Z}_{2}$. Let $\rho: Z_{2} \rightarrow \mathrm{U}(1)$ be the non-trivial representation with $\rho(g)=-1$ where $g$ is the generator. Let $L=S^{3} \times \mathbf{C}$ and identify $(x, z)=(-x,-z)$ to define a line bundle $L_{\rho}$ over $\mathbf{R} P_{3}$ with holonomy $\rho$. We show $L_{\rho}$ is non-trivial. Suppose the contrary, then $L_{\rho}$ is trivial over $\mathbf{R} P_{2}$ as well. This shows there is a map $f: S^{2} \rightarrow S^{1}$ with $-f(-x)=f(x)$. If we restrict $f$ to the upper hemisphere of $S^{2}$, then $f: D_{+}^{2} \rightarrow S^{1}$ satisfies $f(x)=-f(-x)$ on the boundary. Therefore $f$ has odd degree. Since $f$ extends to $D_{+}, f$ must have zero degree. This contradiction establishes no such $f$ exists and $L_{\rho}$ is non-trivial.

The bundle $L_{\rho} \oplus L_{\rho}$ is $S^{3} \times \mathbf{C}$ modulo the relation $(x, z)=(-x,-z)$. We let $g(x): S^{3} \rightarrow \mathrm{SU}(2)$ be the identity map:

$$
g(x)=\left(\begin{array}{cc}
x_{0}+i x_{1} & -x_{2}+i x_{3} \\
x_{2}+i x_{3} & x_{0}-i x_{1}
\end{array}\right)
$$

then $g(-x)=-g(x)$. Thus $g$ descends to give a global frame on $L_{\rho} \oplus L_{\rho}$ so this bundle is topologically trivial and $L_{\rho}$ represents a $\mathbf{Z}_{2}$ torsion class in $K\left(\mathbf{R} P_{3}\right)$. Since $L_{\rho}$ is a line bundle and is not topologically trivial, $L_{\rho}-1$ is a non-zero element of $\widetilde{K}\left(\mathbf{R} P_{3}\right)$. This construction generalizes to define $L_{\rho}$ over $\mathbf{R} P_{n}$. We use the map $g: S^{n} \rightarrow \mathrm{U}\left(2^{k}\right)$ defined by Clifford algebras so $g(x)=-g(-x)$ to show $2^{k} L_{\rho}-2^{k}=0$ in $\widetilde{K}\left(\mathbf{R} P_{n}\right)$ where $k=[n / 2]$; we refer to Lemma 3.8.9 for details.

We suppose henceforth in this section that the bundle $V_{\rho}$ is topologically trivial and let $\vec{s}$ be a global frame for $V_{\rho}$. (In section 4.9 we will study the
more general case). We can take $P_{k}=P \otimes 1$ relative to the frame $\vec{s}$ so that both $P_{k}$ and $P_{\rho}$ are defined on the same bundle with the same leading symbol. We form the 1-parameter family $t P_{k}+(1-t) P_{\rho}=P(t, \rho, \vec{s})$ and define:

$$
\operatorname{ind}(\rho, P, \vec{s})=\int_{0}^{1} \frac{d}{d t} \tilde{\eta}(P(t, \rho, \vec{s})) d t=\int_{M} a(x, \rho, P, \vec{s}) \mathrm{dvol}(x)
$$

where

$$
a(x, \rho, P, \vec{s})=\int_{0}^{1} a\left(x, P(t, \rho, \vec{s}), P^{\prime}(t, \rho, \vec{s})\right) d t
$$

The choice of a global frame permits us to lift ind from $\mathbf{R} \bmod \mathbf{Z}$ to $\mathbf{R}$. If we take the operator and representation of example 4.4.2, then $\operatorname{ind}\left(\rho_{\varepsilon}, P, \vec{s}\right)=\varepsilon$ and thus in particular the lift depends on the global frame chosen (or equivalently on the particular presentation of the representation $\rho$ on a trivial bundle). This also permits us to construct non-trivial real valued invariants even on simply connected manifolds by choosing suitable inequivalent global trivializations of $V_{\rho}$.
Lemma 4.4.4.
(a) $\operatorname{ind}(\rho, P, \vec{s})=\int_{M} a(x, \rho, P, \vec{s}) \operatorname{dvol}(x)$ is given by a local formula which depends on the connection 1-form of $\nabla_{\rho}$ relative to the global frame $\vec{s}$ and on the symbol of the operator $P$.
(b) If $\vec{s}_{t}$ is a smooth 1-parameter family of global sections, then $\operatorname{ind}\left(\rho, P, \vec{s}_{t}\right)$ is independent of the parameter $t$.
Proof: The first assertion follows from the definition of $a(x, \rho, P, \vec{s})$ given above and from the results of the first chapter. This shows $\operatorname{ind}\left(\rho, P, \vec{s}_{t}\right)$ varies continuously with $t$. Since its mod $\mathbf{Z}$ reduction is $\operatorname{ind}(\rho, P)$, this $\bmod \mathbf{Z}$ reduction is constant. This implies $\operatorname{ind}\left(\rho, P, \vec{s}_{t}\right)$ itself is constant.

We can use $\operatorname{ind}(\rho, P, \vec{s})$ to detect inequivalent trivializations of a bundle and thereby study the homotopy $[M, \mathrm{U}(k)]$ even if $M$ is simply connected. This is related to spectral flow.

The secondary characteristic classes are cohomological invariants of the representation $\rho$. They are normally $\mathbf{R} \bmod \mathbf{Z}$ classes, but can be lifted to $\mathbf{R}$ and expressed in terms of local invariants if the bundle $V_{\rho}$ is given a fixed trivialization. We first recall the definition of the Chern character. Let $W$ be a smooth vector bundle with connection $\nabla$. Relative to some local frame, we let $\omega$ be the connection 1-form and $\Omega=d \omega-\omega \wedge \omega$ be the curvature. The Chern character is given by $\operatorname{ch}_{k}(\nabla)=\left(\frac{i}{2 \pi}\right)^{k} \frac{1}{k!} \operatorname{Tr}\left(\Omega^{k}\right)$. This is a closed $2 k$ form independent of the frame $\vec{s}$ chosen. If $\nabla_{i}$ are two connections for $i=0,1$, we form $\nabla_{t}=t \nabla_{1}+(1-t) \nabla_{0}$. If $\theta=\omega_{1}-\omega_{0}$, then $\theta$ transforms like a tensor. If $\Omega_{t}$ is the curvature of the connection $\nabla_{t}$, then:

$$
c h_{k}\left(\nabla_{1}\right)-\operatorname{ch}\left(\nabla_{0}\right)=\int_{0}^{1} \frac{d}{d t} c h_{k}\left(\nabla_{t}\right) d t=d\left(\operatorname{Tch}_{k}\left(\nabla_{1}, \nabla_{0}\right)\right)
$$

where the transgression $T c h_{k}$ is defined by:

$$
\operatorname{Tch}_{k}\left(\nabla_{1}, \nabla_{0}\right)=\left(\frac{i}{2 \pi}\right)^{k} \frac{1}{(k-1)!} \operatorname{Tr}\left\{\int_{0}^{1} \theta \Omega_{t}^{k-1} d t\right\}
$$

We refer to the second chapter for further details on this construction.
We apply this construction to the case in which both $\nabla_{1}$ and $\nabla_{0}$ have zero curvature. We choose a local frame so $\omega_{0}=0$. Then $\omega_{1}=\theta$ and $\Omega_{1}=d \theta-\theta \wedge \theta=0$. Consequently:

$$
\omega_{t}=t \theta \quad \text { and } \quad \Omega_{t}=t d \theta-t^{2} \theta \wedge \theta=\left(t-t^{2}\right) \theta \wedge \theta
$$

so that:

$$
\operatorname{Tch}_{k}\left(\nabla_{1}, \nabla_{0}\right)=\left(\frac{i}{2 \pi}\right)^{k} \frac{1}{(k-1)!} \int_{0}^{1}\left(t-t^{2}\right)^{k-1} d t \cdot \operatorname{Tr}\left(\theta^{2 k-1}\right)
$$

We integrate by parts to evaluate this coefficient:

$$
\begin{aligned}
\int_{0}^{1}\left(t-t^{2}\right)^{k-1} d t & =\int_{0}^{1} t^{k-1}(1-t)^{k-1} d t=\frac{k-1}{k} \int_{0}^{1} t^{k}(1-t)^{k-2} d t \\
& =\frac{(k-1)!}{k \cdot(k+1) \cdots(2 k-2)} \int_{0}^{1} t^{2 k-2} d t \\
& =\frac{(k-1)!(k-1)!}{(2 k-1)!}
\end{aligned}
$$

Therefore

$$
\operatorname{Tch}_{k}\left(\nabla_{1}, \nabla_{0}\right)=\left(\frac{i}{2 \pi}\right)^{k} \frac{(k-1)!}{(2 k-1)!} \cdot \operatorname{Tr}\left(\theta^{2 k-1}\right)
$$

We illustrate the use of secondary characteristic classes by giving another version of the Atiyah-Singer index theorem. Let $Q: C^{\infty}\left(1^{k}\right) \rightarrow C^{\infty}\left(1^{k}\right)$ be an elliptic complex. Let $\vec{s}_{ \pm}$be global frames on $\Pi_{+}(\Sigma q)$ over $D_{ \pm}\left(T^{*} M\right)$ so that $\vec{s}_{-}=q^{t}(x, \xi) \vec{s}_{+}$on $D_{+}\left(T^{*} M\right) \cap D_{-}\left(T^{*} M\right)=S\left(T^{*} M\right)$. (The clutching function is $q$; to agree with the notation adopted in the third section we express $\vec{s}_{-}=q^{t} \vec{s}_{+}$as we think of $q$ being a matrix acting on column vectors of $\mathbf{C}^{k}$. The action on the frame is therefore the transpose action). We choose connections $\nabla_{ \pm}$on $\Pi_{+}(\Sigma q)$ so $\nabla_{ \pm}\left(\vec{s}_{ \pm}\right)=0$ on $D_{ \pm}\left(T^{*} M\right)$. Then:

$$
\begin{aligned}
\operatorname{index}(Q) & =(-1)^{m} \int_{\Sigma\left(T^{*} M\right)} \operatorname{Todd}(M) \wedge \operatorname{ch}\left(\nabla_{-}\right) \\
& =(-1)^{m} \int_{D_{+}(M)} \operatorname{Todd}(M) \wedge \operatorname{ch}\left(\nabla_{-}\right)
\end{aligned}
$$

However, on $D_{+}$we have $\Omega_{+}=0$ so we can replace $\operatorname{ch}\left(\nabla_{-}\right)$by $\operatorname{ch}\left(\nabla_{-}\right)-$ $\operatorname{ch}\left(\nabla_{+}\right)$without changing the value of the integral. $\operatorname{ch}\left(\nabla_{-}\right)-\operatorname{ch}\left(\nabla_{+}\right)=$ $d \operatorname{Tch}\left(\nabla_{-}, \nabla_{+}\right)$so an application of Stokes theorem (together with a careful consideration of the orientations involved) yields:

$$
\operatorname{index}(Q)=(-1)^{m} \int_{S\left(T^{*} M\right)} \operatorname{Todd}(M) \wedge \operatorname{Tch}\left(\nabla_{-}, \nabla_{+}\right)
$$

Both connections have zero curvature near the equator $S\left(T^{*} M\right)$. The fibers over $D_{+}$are glued to the fibers over $D_{-}$using the clutching function $q$. With the notational conventions we have established, if $f_{+}$is a smooth section relative to the frame $\vec{s}_{+}$then the corresponding representation is $q f_{+}$relative to the frame $\vec{s}_{-}$. Therefore $\nabla_{-}\left(f_{+}\right)=q^{-1} d q \cdot f_{+}+d f_{+}$and consequently

$$
\nabla_{-}-\nabla_{+}=\theta=q^{-1} d q
$$

(In obtaining the Maurer-Cartan form one must be careful which convention one uses-right versus left-and we confess to having used both conventions in the course of this book.) We use this to compute:

$$
T c h_{k}\left(\nabla_{-}, \nabla_{+}\right)=\sum c_{k} \operatorname{Tr}\left(\left(q^{-1} d q\right)^{2 k-1}\right)
$$

for

$$
c_{k}=\left(\frac{i}{2 \pi}\right)^{k} \frac{(k-1)!}{(2 k-1)!}
$$

Let $\theta=g^{-1} d g$ be the Maurer-Cartan form and let $T c h=\sum_{k} c_{k} \operatorname{Tr}\left(\theta^{2 k-1}\right)$. This defines an element of the odd cohomology of $\operatorname{GL}(\cdot, \mathbf{C})$ such that $\operatorname{Tch}\left(\nabla_{-}, \nabla_{+}\right)=q^{*}(T c h)$. We summarize these computations as follows:
Lemma 4.4.5. Let $\theta=g^{-1} d g$ be the Maurer-Cartan form on the general linear group. Define:

$$
T c h=\sum_{k}\left(\frac{i}{2 \pi}\right)^{k} \frac{(k-1)!}{(2 k-1)!} \cdot \operatorname{Tr}\left(\theta^{2 k-1}\right)
$$

as an element of the odd cohomology. If $Q: C^{\infty}\left(1^{\cdot}\right) \rightarrow C\left(1^{\cdot}\right)$ is an elliptic complex defined on the trivial bundle, let $q$ be the symbol so $q: S\left(T^{*} M\right) \rightarrow$ $G L(\cdot, \mathbf{C})$. Then $\operatorname{index}(Q)=(-1)^{m} \int_{S\left(T^{*} M\right)} \operatorname{Todd}(M) \wedge q^{*}(T c h)$.

We compute explicitly the first few terms in the expansion:

$$
\operatorname{Tch}=\frac{i}{2 \pi} \operatorname{Tr}(\theta)+\frac{-1}{24 \pi^{2}} \operatorname{Tr}\left(\theta^{3}\right)+\frac{-i}{960 \pi^{3}} \operatorname{Tr}\left(\theta^{5}\right)+\cdots
$$

In this version, the Atiyah-Singer theorem generalizes to the case of $d M \neq \emptyset$ as the Atiyah-Bott theorem. We shall discuss this in section 4.5.

We can now state the Atiyah-Patodi-Singer twisted index theorem:

THEOREM 4.4.6. Let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be an elliptic self-adjoint pseudo-differential operator of order $d>0$. Let $\rho: \pi_{1}(M) \rightarrow \mathrm{U}(k)$ be a unitary representation of the fundamental group and assume the associated bundle $V_{\rho}$ is toplogically trivial. Let $\vec{s}$ be a global frame for $V_{\rho}$ and let $\nabla_{0}(\vec{s}) \equiv 0$ define the connection $\nabla_{0}$. Let $\nabla_{\rho}$ be the connection defined by the representation $\rho$ and let $\theta=\nabla_{\rho}(\vec{s})$ so that:

$$
\operatorname{Tch}\left(\nabla_{\rho}, \nabla_{0}\right)=\sum_{k}\left(\frac{i}{2 \pi}\right)^{k} \frac{(k-1)!}{(2 k-1)!} \cdot \operatorname{Tr}\left(\theta^{2 k-1}\right)
$$

Then:

$$
\operatorname{ind}(\rho, P, \vec{s})=(-1)^{m} \int_{S\left(T^{*} M\right)} \operatorname{Todd}(M) \wedge \operatorname{ch}\left(\Pi_{+} p\right) \wedge \operatorname{Tch}\left(\nabla_{\rho}, \nabla_{0}\right)
$$

$S\left(T^{*} M\right)$ is given the orientation induced by the simplectic orientation $d x_{1} \wedge$ $d \xi_{1} \wedge \cdots \wedge d x_{m} \wedge d \xi_{m}$ on $T^{*} M$ where we use the outward pointing normal so $N \wedge \omega_{2 m-1}=\omega_{2 m}$.

We postpone the proof of Theorem 4.4.6 for the moment to return to the examples considered previously. In example 4.4.2, we let $\rho=e^{2 \pi i \varepsilon}$ on the generator of $\pi_{1} S^{1} \simeq \mathbf{Z}$. We let " 1 " denote the usual trivialization of the bundle $S^{1} \times \mathbf{C}$ so that:

$$
\nabla_{\rho}(1)=-i \varepsilon d \theta \quad \text { and } \quad \operatorname{Tch}\left(\nabla_{\rho}, \nabla_{0}\right)=\varepsilon d \theta / 2 \pi .
$$

The unit sphere bundle decomposes $S\left(T^{*} M\right)=S^{1} \times\{1\} \cup S^{1} \times\{-1\}$. The symbol of the operator $P$ is multiplication by the dual variable $\xi$ so $\operatorname{ch}\left(\Pi_{+} p\right)=1$ on $S^{1} \times\{1\}$ and $\operatorname{ch}\left(\Pi_{+} p\right)=0$ on $S^{1} \times\{-1\}$. The induced orientation on $S^{1} \times\{1\}$ is $-d \theta$. Since $(-1)^{m}=-1$, we compute:

$$
-\int_{S\left(T^{*} M\right)} \operatorname{Todd}(M) \wedge \operatorname{ch}\left(\Pi_{+} p\right) \wedge \operatorname{Tch}\left(\nabla_{\rho}, \nabla_{0}\right)=\int_{0}^{2 \pi} \varepsilon d \theta / 2 \pi=\varepsilon
$$

This also gives an example in which $T c h_{1}$ is non-trivial.
In example 4.4.3, we took the non-trivial representation $\rho$ of $\pi_{1}\left(\mathbf{R} P_{3}\right)=$ $\mathbf{Z}_{2}$ to define a non-trivial complex line bundle $V_{\rho}$ such that $V_{\rho} \oplus V_{\rho}=1^{2}$. The appropriate generalization of this example is to 3-dimensional lens spaces and provides another application of Theorem 4.4.6:
Example 4.4.7: Let $m=3$ and let $n$ and $q$ be relatively prime positive integers. Let $\lambda=e^{2 \pi i / n}$ be a primitive $n^{\text {th }}$ root of unity and let $\gamma=$ $\operatorname{diag}\left(\lambda, \lambda^{q}\right)$ generate a cyclic subgroup $\Gamma$ of $\mathrm{U}(2)$ of order $n$. If $\rho_{s}(\gamma)=\lambda^{s}$, then $\left\{\rho_{s}\right\}_{0 \leq s<n}$ parametrize the irreducible representations of $\Gamma$. $\Gamma$ acts without fixed points on the unit sphere $S^{3}$. Let $L(n, q)=S^{3} / \Gamma$ be the
quotient manifold. As $\pi_{1}\left(S^{3}\right)=0$, we conclude $\pi_{1}(L(n, q))=\Gamma$. Let $V_{s}$ be the line bundle corresponding to the representation $\rho_{s}$. It is defined from $S^{3} \times \mathbf{C}$ by the equivalence relation $\left(z_{1}, z_{2}, w\right)=\left(\lambda z_{1}, \lambda^{q} z_{2}, \lambda^{s} w\right)$.

Exactly the same arguments (working mod $n$ rather than mod 2) used to show $V_{1}$ is non-trivial over $\mathbf{R} P_{3}$ show $V_{s}$ is non-trivial for $0<s<n$. In fact $\widetilde{K}(L(n, q))=\mathbf{Z}_{n}$ as we shall see later in Corollary 4.6.10 and the bundle ( $V_{1}-1$ ) generates the reduced $K$-theory group. A bundle $V_{s_{1}} \oplus \cdots \oplus V_{s_{t}}$ is topologically trivial if and only if $s_{1}+\cdots+s_{t} \equiv 0(n)$.

The bundle $V_{s} \oplus V_{-s}$ is topologically trivial. We define:

$$
g\left(z_{1}, z_{2}\right)=\left(\begin{array}{cc}
z_{1}^{s} & z_{2}^{s q^{\prime}} \\
\bar{z}_{2}^{s q^{\prime}} & -\bar{z}_{1}^{s}
\end{array}\right) \quad \text { where } q q^{\prime} \equiv 1(n)
$$

It is immediate that:

$$
g\left(\lambda z_{1}, \lambda^{q} z_{2}\right)=\left(\begin{array}{cc}
\lambda^{s} z_{1}^{s} & \lambda^{s} z_{2}^{s q^{\prime}} \\
\lambda^{-s} \bar{z}_{2}^{s q^{\prime}} & -\lambda^{-s} \bar{z}_{1}^{s}
\end{array}\right)=\gamma g\left(z_{1}, z_{2}\right)
$$

so we can regard $g$ as an equivariant frame to $V_{s} \otimes V_{-s}$. If $\left\{\nabla_{s} \oplus \nabla_{-s}\right\}$ denotes the connection induced by the locally flat structure and if $\nabla_{0} \oplus \nabla_{0}$ denotes the connection defined by the new frame, then:

$$
\theta=\left\{\nabla_{s} \oplus \nabla_{-s}\right\}-\left\{\nabla_{0} \oplus \nabla_{0}\right\}=d g \cdot g^{-1}
$$

Suppose first $s=q=q^{\prime}=1$ so $g: S^{3} \rightarrow \mathrm{SU}(2)$ is the identity map. $\operatorname{Tr}\left(\theta^{3}\right)$ is a right invariant 3 -form and is therefore a constant multiple of the volume element of $\mathrm{SU}(2)$. We calculate at $z=(1,0)$ in $\mathbf{C}^{2}$. Let

$$
e_{0}=\left(\begin{array}{rr}
1 & 0 \\
0 & -1
\end{array}\right), \quad e_{1}=\left(\begin{array}{rr}
i & 0 \\
0 & i
\end{array}\right), \quad e_{2}=\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right), \quad e_{3}=\left(\begin{array}{rr}
0 & i \\
-i & 0
\end{array}\right)
$$

so that $g(z)=g(x)=x_{0} e_{0}+x_{1} e_{1}+x_{2} e_{2}+x_{3} e_{3}$. At $(1,0)$ we have:

$$
\left(d g \cdot g^{-1}\right)^{3}=\left\{\left(e_{1} d x_{1}+e_{2} d x_{2}+e_{3} d x_{3}\right) e_{0}\right\}^{3} .
$$

$e_{0}$ commutes with $e_{1}$ and anti-commutes with $e_{2}$ and $e_{3}$ so that, as $e_{0}^{2}=1$,

$$
\begin{aligned}
\left(d g \cdot g^{-1}\right)^{3}= & \left(e_{1} d x_{1}+e_{2} d x_{2}+e_{3} d x_{3}\right)\left(e_{1} d x_{1}-e_{2} d x_{2}-e_{3} d x_{3}\right) \times \\
& \left(e_{1} d x_{1}+e_{2} d x_{2}+d e_{3} d x_{3}\right) e_{0} \\
= & \left(-e_{1} e_{2} e_{3}+e_{1} e_{3} e_{2}-e_{2} e_{1} e_{3}-e_{2} e_{3} e_{1}+e_{3} e_{1} e_{2}+e_{3} e_{2} e_{1}\right) \times \\
& e_{0} d x_{1} \wedge d x_{2} \wedge d x_{3} \\
= & -6 e_{1} e_{2} e_{3} e_{0} \cdot \text { dvol }=-6\left(\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right) \text { dvol }
\end{aligned}
$$

so that $\operatorname{Tr}\left(\theta^{3}\right)=-12 \cdot$ dvol. Therefore:

$$
\int_{s^{3}} T c h_{2}\left(\nabla_{1} \oplus \nabla_{-1}, \nabla_{0} \oplus \nabla_{0}\right)=\frac{12}{24 \pi^{2}} \cdot \operatorname{volume}\left(S^{3}\right)=1
$$

We can study other values of $(s, q)$ by composing with the map $\left(z_{1}, z_{2}\right) \mapsto$ $\left(z_{1}^{s}, z_{2}^{s q^{\prime}}\right) / \mid\left(z_{1}^{s}, z_{2}^{s q^{\prime}} \mid\right.$. This gives a map homotopic to the original map and doesn't change the integral. This is an $s^{2} q^{\prime}$-to-one holomorophic map so the corresponding integral becomes $s^{2} q^{\prime}$. If instead of integrating over $S^{3}$ we integrate over $L(n, q)$, we must divide the integral by $n$ so that:

$$
\int_{L(n, q)} T c h_{2}\left(\nabla_{s} \oplus \nabla_{-s}, \nabla_{0} \oplus \nabla_{0}\right)=\frac{s^{2} q^{\prime}}{n}
$$

Let $P$ be the tangential operator of the signature complex. By Lemma 4.3.2, we have:

$$
\int_{s^{2}} \operatorname{ch}\left(\Pi_{+}(P)\right)=-4
$$

$S^{3}$ is parallelizable. The orientation of $S\left(T^{*} S^{3}\right)$ is $d x_{1} \wedge d x_{2} \wedge d \xi_{2} \wedge d x_{3} \wedge$ $d \xi_{3}=-d x_{1} \wedge d x_{2} \wedge d x_{3} \wedge d \xi_{2} \wedge d \xi_{3}$ so $S\left(T^{*} S^{3}\right)=-S^{3} \times S^{2}$ given the usual orientation. Thus

$$
\begin{aligned}
(-1)^{3} \int_{S\left(T^{*} S^{3}\right)} T c h_{2}\left(\nabla_{s} \oplus \nabla_{-s}\right. & \left., \nabla_{0} \oplus \nabla_{0}\right) \wedge c h\left(\Pi_{+} P\right) \\
& =(-1)(-1)(-1) \cdot 4 \cdot s^{2} q^{\prime} / n=-4 s^{2} q^{\prime} / n
\end{aligned}
$$

Consequently by Theorem 4.4.6, we conclude, since $\operatorname{Todd}\left(S^{3}\right)=1$,

$$
\operatorname{ind}\left(\rho_{s}+\rho_{-s}, \rho_{0}+\rho_{0}, P\right)=-4 s^{2} q^{\prime} / n
$$

using the given framing.
The operator $P$ splits into an operator on even and odd forms with equal eta invariants and a corresponding calculation shows that $\operatorname{ind}\left(\rho_{s}+\rho_{-s}\right.$, $\left.\rho_{0}+\rho_{0}, P_{\text {even }}\right)=-2 s^{2} q^{\prime} / n$. There is an orientation preserving isometry $T: L(p, q) \rightarrow L(p, q)$ defined by $T\left(z_{1}, z_{2}\right)=\left(\bar{z}_{1}, \bar{z}_{2}\right)$. It is clear $T$ interchanges the roles of $\rho_{s}$ and $\rho_{-s}$ so that as $\mathbf{R} \bmod \mathbf{Z}$-valued invariants, $\operatorname{ind}\left(\rho_{s}, P_{\text {even }}\right)=\operatorname{ind}\left(\rho_{-s}, P_{\text {even }}\right)$ so that:

$$
\begin{aligned}
\operatorname{ind}\left(\rho_{s}, P\right) & =2 \operatorname{ind}\left(\rho_{s}, P_{\text {even }}\right)=\operatorname{ind}\left(\rho_{s}, P_{\text {even }}\right)+\operatorname{ind}\left(\rho_{-s}, P_{\text {even }}\right) \\
& =-2 s^{2} q^{\prime} / n
\end{aligned}
$$

This gives a formula in $\mathbf{R} \bmod \mathbf{Z}$ for the index of a representation which need not be topologically trivial. We will return to this formula to discuss generalizations in Lemma 4.6.3 when discussing ind $(*, *)$ for general spherical space forms in section 4.6.

We sketch the proof of Theorem 4.4.6. We shall omit many of the details in the interests of brevity. We refer to the papers of Atiyah, Patodi, and Singer for complete details. An elementary proof is contained in (Gilkey, The eta invariant and the secondary characteristic classes of locally flat bundles). Define:

$$
\operatorname{ind}_{1}(\rho, P, s)=(-1)^{m} \int_{S\left(T^{*} M\right)} \operatorname{Todd}(M) \wedge \operatorname{ch}\left(\Pi_{+} p\right) \wedge \operatorname{Tch}\left(\nabla_{\rho}, \nabla_{0}\right)
$$

then Lemmas 4.3.6 and 4.3.7 generalize immediately to:
Lemma 4.4.8.
(a) Let $M_{1}$ and $M_{2}$ be smooth manifolds. Let $P: C^{\infty}\left(V_{1}\right) \rightarrow C^{\infty}\left(V_{2}\right)$ be an elliptic complex over $M_{1}$. Let $Q: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be a self-adjoint elliptic operator over $M_{2}$. We assume $P$ and $Q$ are differential operators of the same order and form $R=\left(\begin{array}{cc}Q & P^{*} \\ P & -Q\end{array}\right)$ over $M=M_{1} \times M_{2}$. Let $\rho$ be a representation of $\pi_{1}\left(M_{2}\right)$. Decompose $\pi_{1}(M)=\pi_{1}\left(M_{1}\right) \oplus \pi_{1}\left(M_{2}\right)$ and extend $\rho$ to act trivially on $\pi_{1}\left(M_{1}\right)$. Then we can identify $V_{\rho}$ over $M$ with the pull-back of $V_{\rho}$ over $M_{2}$. Let $s$ be a global trivialization of $V_{\rho}$ then:

$$
\begin{aligned}
\operatorname{ind}(\rho, R, \vec{s}) & =\operatorname{index}(P) \operatorname{ind}(\rho, Q, \vec{s}) \\
\operatorname{ind}_{1}(\rho, R, \vec{s}) & =\operatorname{index}(P) \operatorname{ind}_{1}(\rho, Q, \vec{s})
\end{aligned}
$$

(b) Let $M_{1}=S^{1}$ be the circle and let $(R, Q)$ be as defined in Lemma 4.3.7, then:

$$
\begin{aligned}
\operatorname{ind}(\rho, R, \vec{s}) & =\operatorname{ind}(\rho, P, \vec{s}) \\
\operatorname{ind}_{1}(\rho, R, \vec{s}) & =\operatorname{ind}_{1}(\rho, P, \vec{s})
\end{aligned}
$$

Proof: The assertions about ind follow directly from Lemmas 4.3.6 and 4.3.7. The assertions about $\mathrm{ind}_{1}$ follow from Lemma 3.9.3(d) and from the Atiyah-Singer index theorem.

Lemma 4.4.8(c) lets us reduce the proof of Theorem 4.4.6 to the case $\operatorname{dim} M$ odd. Since both ind and $\operatorname{ind}_{1}$ are given by local formulas, we may assume without loss of generality that $M$ is also orientable. Using the same arguments as those given in subsection 4.3, we can interpret both ind and $\operatorname{ind}_{1}$ as maps in $K$-theory once the representation $\rho$ and the global frame $\vec{s}$ are fixed. Consequently, the same arguments as those given for the proof of Theorem 4.3.8 permit us to reduce the proof of Theorem 4.4.6 to the case in which $P=A_{V}$ is the operator discussed in Lemma 4.3.4.
$\operatorname{ind}(\rho, P, \vec{s})$ is given by a local formula. If we express everything with respect to the global frame $\vec{s}$, then $P_{k}=P \otimes 1$ and $P_{\rho}$ is functorially expressible in terms of $P$ and in terms of the connection 1-form $\omega=\nabla_{\rho} \vec{s}$;

$$
\operatorname{ind}(\rho, P, \vec{s})=\int_{M} a(x, G, \omega) \operatorname{dvol}(x)
$$

The same arguments as those given in discussing the signature complex show $a(x, G, \omega)$ is homogeneous of order $n$ in the jets of the metric and of the connection 1-form. The local invariant changes sign if the orientation is reversed and thus $a(x, G, \omega)$ dvol $(x)$ should be regarded as an $m$-form not as a measure. The additivity of $\eta$ with respect to direct sums shows $a\left(x, G, \omega_{1} \oplus \omega_{2}\right)=a\left(x, G, \omega_{1}\right)+a\left(x, G, \omega_{2}\right)$ when we take the direct sum of representations. Arguments similar to those given in the second chapter (and which are worked out elsewhere) prove:

Lemma 4.4.9. Let $a(x, G, \omega)$ be an $m$-form valued invariant which is defined on Riemannian metrics and on 1-form valued tensors $\omega$ so that $d \omega-\omega \wedge \omega=0$. Suppose $a$ is homogeneous of order $m$ in the jets of the metric and the tensor $\omega$ and suppose that $a\left(x, G, \omega_{1} \oplus \omega_{2}\right)=a\left(x, G, \omega_{1}\right)+$ $a\left(x, G, \omega_{2}\right)$. Then we can decompose:

$$
a(x, G, \omega)=\sum_{\nu} f_{\nu}(G) \wedge \operatorname{Tch}_{\nu}(\omega)=\sum_{\nu} f_{\nu}(G) \wedge c_{\nu} \operatorname{Tr}\left(\omega^{2 \nu-1}\right)
$$

where $f_{\nu}(G)$ are real characteristic forms of $T(M)$ of order $m+1-2 \nu$.
We use the same argument as that given in the proof of the Atiyah-Singer index theorem to show there must exist a local formula for $\operatorname{ind}(\rho, P, s)$ which has the form:

$$
\begin{aligned}
& \operatorname{ind}(\rho, P, s)= \\
& \quad \int_{S\left(T^{*} M\right)} \sum_{4 i+2 j+2 k=m+1} T d_{i, m}^{\prime}(M) \wedge c h_{j}\left(\Pi_{+} p\right) \wedge \operatorname{Tch}_{k}\left(\nabla_{\rho}, \nabla_{0}\right)
\end{aligned}
$$

When the existence of such a formula is coupled with the product formula given in Lemma 4.4.8(a) and with the Atiyah-Singer index theorem, we deduce that the formula must actually have the form:

$$
\begin{aligned}
& \operatorname{ind}(\rho, P, s)= \\
& \quad \int_{S\left(T^{*} M\right)} \sum_{j+2 k=m+1}\left(\operatorname{Todd}(M) \wedge c h\left(\Pi_{+} p\right)\right)_{j} \wedge c(k) \operatorname{Tch}_{k}\left(\nabla_{\rho}, \nabla_{0}\right)
\end{aligned}
$$

where $c(k)$ is some universal constant which remains to be determined.
If we take an Abelian representation, all the $T c h_{k}$ vanish for $k>1$. We already verified that the constant $c(1)=1$ by checking the operator of example 4.4.2 on the circle. The fact that the other normalizing constants are also 1 follows from a detailed consideration of the asymptotics of the heat equation which arise; we refer to (Gilkey, The eta invariant and the secondary characteristic classes of locally flat bundles) for further details regarding this verification. Alternatively, the Atiyah-Patodi-Singer index
theorem for manifolds with boundary can be used to check these normalizing constants; we refer to (Atiyah, Patodi, Singer: Spectral asymmetry and Riemannian geometry I-III) for details.

The Atiyah-Singer index theorem is a formula on $K\left(\Sigma\left(T^{*} M\right)\right)$. The Atiyah-Patodi-Singer twisted index theorem can be regarded as a formula on $K\left(S\left(T^{*} M\right)\right)=K^{1}\left(\Sigma\left(T^{*} M\right)\right)$. These two formulas are to be regarded as suspensions of each other and are linked by Bott periodicity in a purely formal sense which we shall not make explicit.

If $V_{\rho}$ is not topologically trivial, there is no local formula for $\operatorname{ind}(\rho, P)$ in general. It is possible to calculate this using the Lefschetz fixed point formulas as we will discuss later.

We conclude this section by discussing the generalization of Theorem 4.4.6 to the case of manifolds with boundary. It is a fairly straightforward computation using the methods of section 3.9 to show:
Lemma 4.4.9. We adopt the notation of Theorem 4.4.6. Then for any $k \geq 0$,
$\operatorname{ind}(\rho, P, \vec{s})=(-1)^{m} \int_{\Sigma^{2 k}\left(T^{*} M\right)} \operatorname{Todd}(M) \wedge \operatorname{ch}\left(\Pi_{+}\left(\Sigma^{2 k}(p)\right)\right) \wedge \operatorname{Tch}\left(\nabla_{\rho}, \nabla_{0}\right)$.
Remark: This shows that we can stabilize by suspending as often as we please. The index can be computed as an integral over $\Sigma^{2 k+1}\left(T^{*} M\right)$ while the twisted index is an integral over $\Sigma^{2 k}\left(T^{*} M\right)$. These two formulas are at least formally speaking the suspensions of each other.

It turns out that the formula of Theorem 4.4.6 does not generalize directly to the case of manifolds with boundary, while the formula of Lemma 4.4.9 with $k=1$ does generalize. Let $M$ be a compact manifold with boundary $d M$. We choose a Riemannian metric on $M$ which is product near $d M$. Let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be a first order partial differential operator with leading symbol $p$ which is formally self-adjoint. We suppose $p(x, \xi)^{2}=|\xi|^{2} \cdot 1_{V}$ so that $p$ is defined by Clifford matrices. If we decompose $p(x, \xi)=\sum_{j} p_{j}(x) \xi_{j}$ relative to a local orthonormal frame for $T^{*}(M)$, then the $p_{j}$ are self-adjoint and satisfy the relations $p_{j} p_{k}+p_{k} p_{j}=2 \delta_{j k}$.

Near the boundary we decompose $T^{*}(M)=T^{*}(d M) \oplus 1$ into tangential and normal directions. We let $\xi=(\zeta, z)$ for $\zeta \in T^{*}(d M)$ reflect this decomposition. Decompose $p(x, \xi)=\sum_{1 \leq j \leq m-1} p_{j}(x) \xi_{j}+p_{m} z$. Let $t$ be a real parameter and define:

$$
\tau(x, \zeta, t)=i \cdot p_{m}\left\{\sum_{1 \leq j \leq m-1} p_{j}(x) \zeta_{j}-i t\right\}
$$

as the endomorphism defined in Chapter 1. We suppose given a self-adjoint endomorphism $q$ of $V$ which anti-commutes with $\tau$. Let $B$ denote the orthogonal projection $\frac{1}{2}(1+q)$ on the +1 eigenspace of $q . \quad(P, B)$ is a selfadjoint elliptic boundary value problem. The results proved for manifolds without boundary extend to this case to become:

Theorem 4.4.10. Let $(P, B)$ be an elliptic first order boundary value problem. We assume $P$ is self-adjoint and $\sigma_{L}(p)^{2}=|\xi|^{2} \cdot I_{V}$. We assume $B=\frac{1}{2}(1+q)$ where $q$ anti-commutes with $\tau=i \cdot p_{m}\left(\sum_{1 \leq j \leq m-1} p_{j} \zeta_{j}-i t\right)$. Let $\left\{\lambda_{\nu}\right\}_{\nu=1}^{\infty}$ denote the spectrum of the operator $P_{B}$. Define:

$$
\eta(s, P, B)=\sum_{\nu} \operatorname{sign}\left(\lambda_{\nu}\right)\left|\lambda_{\nu}\right|^{-s}
$$

(a) $\eta(s, P, B)$ is well defined and holomorphic for $\operatorname{Re}(s) \gg 0$.
(b) $\eta(s, P, B)$ admits a meromorphic extension to $\mathbf{C}$ with isolated simple poles at $s=(n-m) / 2$ for $n=0,1,2, \ldots$. The residue of $\eta$ at such a simple pole is given by integrating a local formula $a_{n}(x, P)$ over $M$ and a local formula $a_{n-1}(x, P, B)$ over $d M$.
(c) The value $s=0$ is a regular value.
(d) If $\left(P_{u}, B_{u}\right)$ is a smooth 1-parameter family of such operators, then the derivative $\frac{d}{d u} \eta\left(0, P_{u}, B_{u}\right)$ is given by a local formula.
Remark: This theorem holds in much greater generality than we are stating it; we restrict to operators and boundary conditions given by Clifford matrices to simplify the discussion. The reader should consult (Gilkey-Smith) for details on the general case.

Such boundary conditions always exist if $M$ is orientable and $m$ is even; they may not exist if $m$ is odd. For $m$ even and $M$ orientable, we can take $q=p_{1} \ldots p_{m-1}$. If $m$ is odd, the obstruction is $\operatorname{Tr}\left(p_{1} \ldots p_{m}\right)$ as discussed earlier. This theorem permits us to define $\operatorname{ind}(\rho, P, B)$ if $\rho$ is a unitary representation of $\pi_{1}(M)$ just as in the case where $d M$ is empty. In $\mathbf{R} \bmod \mathbf{Z}$, it is a homotopy invariant of $(P, B)$. If the bundle $V_{\rho}$ is topologically trivial, we can define

$$
\operatorname{ind}(\rho, P, B, \vec{s})
$$

as a real-valued invariant which is given by a local formula integrated over $M$ and $d M$.

The boundary condition can be used to define a homotopy of $\Sigma p$ to an elliptic symbol which doesn't depend upon the tangential fiber coordinates. Homotopy 4.4.11: Let $(P, B)$ be as in Theorem 4.4.10 with symbols given by Clifford matrices. Let $u$ be an auxilary parameter and define:

$$
\tau(x, \zeta, t, u)=\cos \left(\frac{\pi}{2} \cdot u\right) \tau(x, \zeta, t)+\sin \left(\frac{\pi}{2} \cdot u\right) q \quad \text { for } u \in[-1,0]
$$

It is immediate that $\tau(x, \zeta, t, 0)=\tau(x, \zeta, t)$ and $\tau(x, \zeta, t,-1)=-q$. As $\tau$ and $q$ are self-adjoint and anti-commute, this is a homotopy through self-adjoint matrices with eigenvalues $\pm\left\{\cos ^{2}\left(\frac{\pi}{2} \cdot u\right)\left(|\zeta|^{2}+t^{2}\right)+\sin ^{2}\left(\frac{\pi}{2} \cdot u\right)\right\}^{1 / 2}$.

Therefore $\tau(x, \zeta, t, u)-i z$ is a non-singular elliptic symbol for $(\zeta, t, z) \neq$ $(0,0,0)$. We have

$$
\Sigma p(x, \zeta, z, t)=p(x, \zeta, z)-i t=i p_{m} \cdot\{\tau(x, \zeta, t)-i z\}
$$

so we define the homotopy

$$
\Sigma p(x, \zeta, z, t, u)=i p_{m} \cdot\{\tau(x, \zeta, t, u)-i z\} \quad \text { for } u \in[-1,0] .
$$

We identifiy a neighborhood of $d M$ in $M$ with $d M \times[0, \delta)$. We sew on a collared neighborhood $d M \times[-1,0]$ to define $\widetilde{M}$. We use the homotopy just defined to define an elliptic symbol $(\Sigma p)_{B}$ on $\Sigma\left(T^{*} \widetilde{M}\right)$ which does not depend upon the tangential fiber coordinates $\zeta$ on the boundary $d M \times$ $\{-1\}$. We use the collaring to construct a diffeomorophism of $M$ and $\widetilde{M}$ to regard $(\Sigma p)_{B}$ on $\Sigma\left(T^{*} M\right)$ where this elliptic symbol is independent of the tangential fiber coordinates on the boundary. We call this construction Homotopy 4.4.11.

We can now state the generalization of Theorem 4.4.6 to manifolds with boundary:

Theorem 4.4.12. We adopt the notation of Theorem 4.4.10 and let $(P, B)$ be an elliptic self-adjoint first order boundary value problem with the symbols given by Clifford matrices. Let $\rho$ be a representation of $\pi_{1}(M)$. Suppose $V_{\rho}$ is topologically trivial and let $\vec{s}$ be a global frame. Then

$$
\begin{aligned}
& \operatorname{ind}(\rho, P, B, \vec{s})= \\
& (-1)^{m} \int_{\Sigma^{2 k}\left(T^{*} M\right)} \operatorname{Todd}(M) \wedge \operatorname{ch}\left(\Pi_{+}\left(\Sigma^{2 k-1}\left((\Sigma p)_{B}\right)\right)\right) \wedge \operatorname{Tch}\left(\nabla_{\rho}, \nabla_{0}\right)
\end{aligned}
$$

for any $k \geq 1$. Here $(\Sigma p)_{B}$ is the symbol on $\Sigma\left(T^{*} M\right)$ defined by Homotopy 4.4.11 so that it is an elliptic symbol independent of the tangential fiber variables $\zeta$ on the boundary.
Remark: This theorem is in fact true in much greater generality. It is true under the much weaker assumption that $P$ is a first order formally self-adjoint elliptic differential operator and that $(P, B)$ is self-adjoint and strongly elliptic in the sense discussed in Chapter 1. The relevant homotopy is more complicated to discuss and we refer to (Gilkey-Smith) for both details on this generalization and also for the proof of this theorem.

We conclude by stating the Atiyah-Bott formula in this framework. Let $P: C^{\infty}\left(V_{1}\right) \rightarrow C^{\infty}\left(V_{2}\right)$ be a first order elliptic operator and let $B$ be an elliptic boundary value problem. The boundary value problem gives a homotopy of $\Sigma p$ through self-adjoint elliptic symbols to a symbol independent of the tangential fiber variables. We call the new symbol $\{\Sigma p\}_{B}$. The Atiyah-Bott formula is

$$
\operatorname{index}(P, B)=(-1)^{m} \int_{\Sigma\left(T^{*} M\right)} \operatorname{Todd}(M) \wedge \operatorname{ch}\left(\Pi_{+}\left((\Sigma p)_{B}\right)\right)
$$

### 4.5. Lefschetz Fixed Point Formulas.

In section 1.8, we discussed the Lefschetz fixed point formulas using heat equation methods. In this section, we shall derive the classical Lefschetz fixed point formulas for a non-degenerate smooth map with isolated fixed points for the four classical elliptic complexes. We shall also discuss the case of higher dimensional fixed point sets for the de Rham complex. The corresponding analysis for the signature and spin complexes is much more difficult and is beyond the scope of this book, and we refer to (Gilkey, Lefschetz fixed point formulas and the heat equation) for further details. We will conclude by discussing the theorem of Donnelly relating Lefschetz fixed point formulas to the eta invariant.

A number of authors have worked on proving these formulas. Kotake in 1969 discussed the case of isolated fixed points. In 1975, Lee extended the results of Seeley to yield the results of section 1.8 giving a heat equation approach to the Lefschetz fixed point formulas in general. We also derived these results independently not being aware of Lee's work. Donnelly in 1976 derived some of the results concerning the existence of the asymptotic expansion if the map concerned was an isometry. In 1978 he extended his results to manifolds with boundary.

During the period 1970 to 1976, Patodi had been working on generalizing his results concerning the index theorem to Lefschetz fixed point formulas, but his illness and untimely death in December 1976 prevented him from publishing the details of his work on the $G$-signature theorem. Donnelly completed Patodi's work and joint papers by Patodi and Donnelly contain these results. In 1976, Kawasaki gave a proof of the $G$-signature theorem in his thesis on $V$-manifolds. We also derived all the results of this section independently at the same time. In a sense, the Lefschetz fixed point formulas should have been derived by heat equation methods at the same time as the index theorem was proved by heat equation methods in 1972 and it remains a historical accident that this was not done. The problem was long over-due for solution and it is not surprising that it was solved simultaneously by a number of people.

We first assume $T: M \rightarrow M$ is an isometry.
Lemma 4.5.1. If $T$ is an isometry, then the fixed point set of $T$ consists of the disjoint union of a finite number of totally geodesic submanifolds $N_{1}, \ldots$. If $N$ is one component of the fixed point set, the normal bundle $\nu$ is the orthogonal complement of $T(N)$ in $\left.T(M)\right|_{N} . \quad \nu$ is invariant under $d T$ and $\operatorname{det}\left(I-d T_{\nu}\right)>0$ so $T$ is non-degenerate.

Proof: Let $a>0$ be the injectivity radius of $M$ so that if $\operatorname{dist}(x, y)<a$, then there exists a unique shortest geodesic $\gamma$ joining $x$ to $y$ in $M$. If $T(x)=x$ and $T(y)=y$, then $T \gamma$ is another shortest geodesic joining $x$ to $y$ so $T \gamma=\gamma$ is fixed pointwise and $\gamma$ is contained in the fixed point
set. This shows the fixed point set is totally geodesic. Fix $T(x)=x$ and decompose $T(M)_{x}=V_{1} \oplus \nu$ where $V_{1}=\left\{v \in T(M)_{x}: d T(x) v=v\right\}$. Since $d T(x)$ is orthogonal, both $V_{1}$ and $\nu$ are invariant subspaces and $d T_{\nu}$ is an orthogonal matrix with no eigenvalue 1. Therefore $\operatorname{det}\left(I-d T_{\nu}\right)>0$. Let $\gamma$ be a geodesic starting at $x$ so $\gamma^{\prime}(0) \in V_{1}$. Then $T \gamma$ is a geodesic starting at $x$ with $T \gamma^{\prime}(0)=d T(x) \gamma^{\prime}(0)=\gamma^{\prime}(0)$ so $T \gamma=\gamma$ pointwise and $\exp _{x}\left(V_{1}\right)$ parametrizes the fixed point set near $x$. This completes the proof.

We let $M$ be oriented and of even dimension $m=2 n$. Let $(d+\delta): C^{\infty}\left(\Lambda^{+}\right)$ $\rightarrow C^{\infty}\left(\Lambda^{-}\right)$be the operator of the signature complex. We let $H^{ \pm}(M ; \mathbf{C})=$ $\mathrm{N}\left(\Delta_{ \pm}\right)$on $C^{\infty}\left(\Lambda^{ \pm}\right)$so that $\operatorname{signature}(M)=\operatorname{dim} H^{+}-\operatorname{dim} H^{-}$. If $T$ is an orientation preserving isometry, then $T^{*} d=d T^{*}$ and $T^{*} *=* T^{*}$ where "*" denotes the Hodge operator. Therefore $T^{*}$ induces maps $T^{ \pm}$on $H^{ \pm}(M ; \mathbf{C})$. We define:

$$
L(T)_{\text {signature }}=\operatorname{Tr}\left(T^{+}\right)-\operatorname{Tr}\left(T^{-}\right)
$$

Let $A=d T \in \operatorname{SO}(m)$ at an isolated fixed point. Define

$$
\operatorname{defect}(A, \operatorname{signature})=\left\{\operatorname{Tr}\left(\Lambda^{+}(A)\right)-\operatorname{Tr}\left(\Lambda^{-}(A)\right)\right\} / \operatorname{det}(I-A)
$$

as the contribution from Lemma 1.8.3. We wish to calculate this characteristic polynomial. If $m=2$, let $\left\{e_{1}, e_{2}\right\}$ be an oriented orthonormal basis for $T(M)=T^{*}(M)$ such that

$$
A e_{1}=(\cos \theta) e_{1}+(\sin \theta) e_{2} \quad \text { and } \quad A e_{2}=(\cos \theta) e_{2}-(\sin \theta) e_{1}
$$

The representation spaces are defined by:
$\Lambda^{+}=\operatorname{span}\left\{1+i e_{1} e_{2}, e_{1}+i e_{2}\right\} \quad$ and $\quad \Lambda^{-}=\operatorname{span}\left\{1-i e_{1} e_{2}, e_{1}-i e_{2}\right\}$
so that:
$\operatorname{Tr}\left(\Lambda^{+}(A)\right)=1+e^{-i \theta}, \quad \operatorname{Tr}\left(\Lambda^{-}(A)\right)=1-e^{-i \theta}, \quad \operatorname{det}(I-A)=2-2 \cos \theta$
and consequently:

$$
\operatorname{defect}(A, \operatorname{signature})=(-2 i \cos \theta) /(2-2 \cos \theta)=-i \cot (\theta / 2)
$$

More generally let $m=2 n$ and decompose $A=d T$ into a product of mutually orthogonal and commuting rotations through angles $\theta_{j}$ corresponding to complex eigenvalues $\lambda_{j}=e^{i \theta_{j}}$ for $1 \leq j \leq n$. The multiplicative nature of the signature complex then yields the defect formula:

$$
\operatorname{defect}(A, \text { signature })=\prod_{j=1}^{n}\left\{-i \cot \left(\theta_{j} / 2\right)\right\}=\prod_{j=1}^{n} \frac{\lambda_{j}+1}{\lambda_{j}-1}
$$

This is well defined since the condition that the fixed point be isolated is just $0<\theta_{j}<2 \pi$ or equivalently $\lambda_{j} \neq 1$.

There are similar characteristic polynomials for the other classical elliptic complexes. For the de Rham complex, we noted in Chapter 1 the corresponding contribution to be $\pm 1=\operatorname{sign}(\operatorname{det}(I-A))$. If $M$ is spin, let $A: C^{\infty}\left(\Delta^{+}\right) \rightarrow C^{\infty}\left(\Delta^{-}\right)$be the spin complex. If $T$ is an isometry which can be lifted to a spin isometry, we can define $L(T)_{\text {spin }}$ to be the Lefschetz number of $T$ relative to the spin complex. This lifts $A$ from $\mathrm{SO}(m)$ to $\operatorname{SPIN}(m)$. We define:

$$
\operatorname{defect}(A, \operatorname{spin})=\left\{\operatorname{Tr}\left(\Delta^{+}(A)-\operatorname{Tr}\left(\Delta^{-}(A)\right)\right\} / \operatorname{det}(I-A)\right.
$$

and use Lemma 3.2.5 to calculate if $m=2$ that:

$$
\begin{aligned}
\operatorname{defect}(A, \text { spin }) & =\left(e^{-i \theta / 2}-e^{i \theta / 2}\right) /(2-2 \cos \theta) \\
& =-i \sin (\theta / 2) /(1-\cos \theta)=-\frac{i}{2} \operatorname{cosec}(\theta / 2) \\
& =\{\sqrt{\bar{\lambda}}-\sqrt{\lambda}\} /(2-\lambda-\bar{\lambda})=\sqrt{\lambda} /(\lambda-1)
\end{aligned}
$$

Using the multiplicative nature we get a similar product formula in general.
Finally, let $M$ be a holomorphic manifold and let $T: M \rightarrow M$ be a holomorphic map. Then $T$ and $\bar{\partial}$ commute. We let $L(T)_{\text {Dolbeault }}=$ $\sum_{q}(-1)^{q} \operatorname{Tr}\left(T^{*}\right.$ on $\left.H^{0, q}\right)$ be the Lefschetz number of the Dolbeault complex. Just because $T$ has isolated fixed points does not imply that it is non-degenerate; the map $z \mapsto z+1$ defines a map on the Riemann sphere $S^{2}$ which has a single isolated degenerate fixed point at $\infty$. We suppose $A=L(T) \in \mathrm{U}\left(\frac{m}{2}\right)$ is in fact non-degenerate and define

$$
\operatorname{defect}(A, \text { Dolbeault })=\left\{\operatorname{Tr}\left(\Lambda^{0, \text { even }}(A)\right)-\operatorname{Tr}\left(\Lambda^{0, \text { odd }}(A)\right)\right\} \operatorname{det}\left(I-A_{\text {real }}\right)
$$

If $m=2$, it is easy to calculate

$$
\begin{aligned}
\operatorname{defect}(A, \text { Dolbeault }) & =\left(1-e^{i \theta}\right) /(2-2 \cos \theta)=(1-\lambda) /(2-\lambda-\bar{\lambda}) \\
& =\lambda /(\lambda-1)
\end{aligned}
$$

with a similar multiplicative formula for $m>2$. We combine these result with Lemma 1.8.3 to derive the classical Lefschetz fixed point formulas:

Theorem 4.5.2. Let $T: M \rightarrow M$ be a non-degenerate smooth map with isolated fixed points at $F(T)=\left\{x_{1}, \ldots, x_{r}\right\}$. Then
(a) $L(T)_{\text {de Rham }}=\sum_{j} \operatorname{sign}(\operatorname{det}(I-d T))\left(x_{j}\right)$.
(b) Suppose $T$ is an orientation preserving isometry. Define:

$$
\operatorname{defect}(A, \text { signature })=\prod_{j}\left\{-i \cdot \cot \theta_{j}\right\}=\prod_{j} \frac{\lambda_{j}+1}{\lambda_{j}-1},
$$

then $L(T)_{\text {signature }}=\sum_{j} \operatorname{defect}\left(d T\left(x_{j}\right)\right.$, signature $)$.
(c) Suppose $T$ is an isometry preserving a spin structure. Define:

$$
\operatorname{defect}(A, \operatorname{spin})=\prod_{j}\left\{-\frac{i}{2} \operatorname{cosec}\left(\theta_{j} / 2\right)\right\}=\prod_{j} \frac{\sqrt{\lambda_{j}}}{\lambda_{j}-1},
$$

then $L(T)_{\text {spin }}=\sum_{j} \operatorname{defect}\left(d T\left(x_{j}\right)\right.$, spin $)$.
(d) Suppose $T$ is holomorphic. Define:

$$
\operatorname{defect}(A, \text { Dolbeault })=\prod_{j} \frac{\lambda_{j}}{\lambda_{j}-1},
$$

then $L(T)_{\text {Dolbeault }}=\sum_{j} \operatorname{defect}\left(d T\left(x_{j}\right)\right.$, Dolbeault $)$.
We proved (a) in Chapter 1. (b)-(d) follow from the calculations we have just given together with Lemma 1.8.3. In defining the defect, the rotation angles of $A$ are $\left\{\theta_{j}\right\}$ so $A e_{1}=(\cos \theta) e_{1}+(\sin \theta) e_{2}$ and $A e_{2}=(-\sin \theta) e_{1}+$ $(\cos \theta) e_{2}$. The corresponding complex eigenvalues are $\left\{\lambda_{j}=e^{i \theta_{j}}\right\}$ where $1 \leq j \leq m / 2$. The formula in (a) does not depend on the orientation. The formula in (b) depends on the orientation, but not on a unitary structure. The formula in (c) depends on the particular lift of $A$ to spin. The formula in (d) depends on the unitary structure. The formulas (b)-(d) are all for even dimensional manifolds while (a) does not depend on the parity of the dimension. There is an elliptic complex called the $\mathrm{PIN}_{c}$ complex defined over non-orientable odd dimensional manifolds which also has an interesting Lefschetz number relative to an orientation reversing isometry.

Using Lemmas 1.8.1 and 1.8.2 it is possible to get a local formula for the Lefschetz number concentrated on the fixed pont set even if $T$ has higher dimensional fixed point sets. A careful analysis of this situation leads to the $G$-signature theorem in full generality. We refer to the appropriate papers of (Gilkey, Kawasaki, and Donnelly) for details regarding the signature and spin complexes. We shall discuss the case of the de Rham complex in some detail.

The interesting thing about the Dolbeault complex is that Theorem 4.5.2 does not generalize to yield a corresponding formula in the case of higher dimensional fixed point sets in terms of characteristic classes. So far, it has not proven possible to identify the invariants of the heat equation with generalized cohomology classes in this case. The Atiyah-Singer index theorem in its full generality also does not yield such a formula. If one assumes that $T$ is an isometry of a Kaehler metric, then the desired result follows by passing first to the $\mathrm{SPIN}_{c}$ complex as was done in Chapter 3. However, in the general case, no heat equation proof of a suitable generalization is yet known. We remark that Toledo and Tong do have a formula
for $L(T)_{\text {Dolbeault }}$ in this case in terms of characteristic classes, but their method of proof is quite different.

Before proceeding to discuss the de Rham complex in some detail, we pause to give another example:
Example 4.5.3: Let $T_{2}$ be the 2-dimensional torus $S^{1} \times S^{1}$ with usual periodic parameters $0 \leq x, y \leq 1$. Let $T(x, y)=(-y, x)$ be a rotation through $90^{\circ}$. Then
(a) $L(T)_{\text {de Rham }}=2$.
(b) $L(T)_{\text {signature }}=2 i$.
(c) $L(T)_{\text {spin }}=2 i / \sqrt{2}$.
(d) $L(T)_{\text {Dolbeault }}=1-i$.

Proof: We use Theorem 4.5.2 and notice that there are two fixed points at $(0,0)$ and $(1 / 2,1 / 2)$. We could also compute directly using the indicated action on the cohomology groups. This shows that although the signature of a manifold is always zero if $m \equiv 2$ (4), there do exist nontrivial $L(T)_{\text {signature }}$ in these dimensions. Of course, there are many other examples.

We now consider the Lefschetz fixed point formula for the de Rham complex if the fixed point set is higher dimensional. Let $T: M \rightarrow M$ be smooth and non-degenerate. We assume for the sake of simplicity for the moment that the fixed point set of $T$ has only a single component $N$ of dimension $n$. The general case will be derived by summing over the components of the fixed point set. Let $\nu$ be the sub-bundle of $\left.T(M)\right|_{N}$ spanned by the generalized eigenvectors of $d T$ corresponding to eigenvalues other than 1. We choose the metric on $M$ so the decomposition $\left.T(M)\right|_{N}=T(N) \oplus \nu$ is orthogonal. We further normalize the choice of metric by assuming that $N$ is a totally goedesic submanifold.

Let $a_{k}\left(x, T, \Delta_{p}\right)$ denote the invariants of the heat equation discussed in Lemma 1.8.1 so that

$$
\operatorname{Tr}\left(T^{*} e^{-t \Delta p}\right) \sim \sum_{k=0}^{\infty} t^{(k-n) / 2} \int_{N} a_{k}\left(x, T, \Delta_{p}\right) \operatorname{dvol}(x)
$$

Let $a_{k}(x, T)_{\text {de Rham }}=\sum(-1)^{p} a_{k}\left(x, T, \Delta_{p}\right)$, then Lemma 1.8.2 implies:

$$
L(T)_{\mathrm{de} \mathrm{Rham}} \sim \sum_{k=0}^{\infty} t^{(k-n) / 2} \int_{N} a_{k}(x, T)_{\mathrm{de} \mathrm{Rham}} \operatorname{dvol}(x)
$$

Thus:

$$
\int_{N} a_{k}(x, T)_{\text {de Rham }} \operatorname{dvol}(x)= \begin{cases}0 & \text { if } k \neq \operatorname{dim} N \\ L(T)_{\text {de Rham }} & \text { if } k=\operatorname{dim} N\end{cases}
$$

We choose coordinates $x=\left(x_{1}, \ldots, x_{n}\right)$ on $N$. We extend these coordinates to a system of coordinates $z=(x, y)$ for $M$ near $N$ where $y=\left(y_{1}, \ldots, y_{m-n}\right)$. We adopt the notational convention:

$$
\begin{aligned}
\text { indices } 1 \leq i, j, k \leq m & \text { index a frame for }\left.T(M)\right|_{N} \\
1 \leq a, b, c \leq n & \text { index a frame for } T(N) \\
n<u, v, w \leq m & \text { index a frame for } \nu .
\end{aligned}
$$

If $\alpha$ is a multi-index, we deompose $\alpha=\left(\alpha_{N}, \alpha_{\nu}\right)$ into tangential and normal components.

We let $g_{i j / \alpha}$ denote the jets of the metric tensor. Let $T=\left(T_{1}, \ldots, T_{m}\right)$ denote the components of the map $T$ relative to the coordinate system $Z$. Let $T_{i / \beta}$ denote the jets of the map $T$; the $T_{i / j}$ variables are tensorial and are just the components of the Jacobian. Define:

$$
\operatorname{ord}\left(g_{i j / \alpha}\right)=|\alpha|, \quad \operatorname{ord}\left(T_{i / \beta}\right)=|\beta|-1, \quad \operatorname{deg}_{v}\left(T_{i / \beta}\right)=\delta_{i, v}+\beta(v)
$$

Let $\mathcal{T}$ be the polynomial algebra in the formal variables $\left\{g_{i j / \alpha}, T_{i / \beta}\right\}$ for $|\alpha|>0,|\beta|>1$ with coefficients $c\left(g_{i j}, T_{i / j}\right)\left|\operatorname{det}\left(I-d T_{\nu}\right)\right|^{-1}$ where the $c\left(g_{i j}, T_{i / j}\right)$ are smooth in the $g_{i j}$ and $T_{i / j}$ variables. The results of section 1.8 imply $a_{k}(x, T)_{\text {de }}$ Rham $\in \mathcal{T}$. If $Z$ is a coordinate system, if $G$ is a metric, and if $T$ is the germ of a non-degenerate smooth map we can evaluate $p(Z, T, G)(x)$ for $p \in \mathcal{T}$ and $x \in N . \quad p$ is said to be invariant if $p(Z, T, G)(x)=p\left(Z^{\prime}, T, G\right)(x)$ for any two coordinate systems $Z, Z^{\prime}$ of this form. We let $\mathcal{T}_{m, n}$ be the sub-algebra of $\mathcal{T}$ of all invariant polynomials and we let $\mathcal{T}_{m, n, k}$ be the sub-space of all invariant polynomials which are homogeneous of degree $k$ using the grading defined above. It is not difficult to show there is a direct sum decomposition $\mathcal{T}_{m, n}=\bigoplus_{k} \mathcal{T}_{m, n, k}$ as a graded algebra.

Lemma 4.5.4. $a_{k}(x, T)_{\text {de Rham }} \in \mathcal{T}_{m, n, k}$.
Proof: The polynomial dependence upon the jets involved together with the form of the coefficients follows from Lemma 1.8.1. Since $a_{k}(x, T)_{\text {de Rham }}$ does not depend upon the coordinate system chosen, it is invariant. We check the homogeneity using dimensional analysis as per usual. If we replace the metric by a new metric $c^{2} G$, then $a_{k}$ becomes $c^{-k} a_{k}$. We replace the coordinate system $Z$ by $Z^{\prime}=c Z$ to replace $g_{i j / \alpha}$ by $c^{-|\alpha|} g_{i j / \alpha}$ and $T_{i / \beta}$ by $c^{-|\beta|+1} T_{i / \beta}$. This completes the proof. The only feature different from the analysis of section 2.4 is the transformation rule for the variables $T_{i / \beta}$.

The invariants $a_{k}$ are multiplicative.
LEMMA 4.5.5. Let $T^{\prime}: M^{\prime} \rightarrow M^{\prime}$ be non-degenerate with fixed submanifold $N^{\prime}$. Define $M=S^{1} \times M^{\prime}, N=S^{1} \times N^{\prime}$, and $T=I \times T^{\prime}$. Then
$T: M \rightarrow M$ is non-degenerate with fixed submanifold $N$. We give $M$ the product metric. Then $a_{k}(x, T)_{\text {de Rham }}=0$.

Proof: We may decompose
$\Lambda^{p}(M)=\Lambda^{0}\left(S^{1}\right) \otimes \Lambda^{p}\left(M^{\prime}\right) \oplus \Lambda^{1}\left(S^{1}\right) \otimes \Lambda^{p-1}\left(M^{\prime}\right) \simeq \Lambda^{p}\left(M^{\prime}\right) \oplus \Lambda^{p-1}\left(M^{\prime}\right)$.
This decomposition is preserved by the map $T$. Under this decomposition, the Laplacian $\Delta_{p}^{M}$ splits into $\Delta_{p}^{\prime} \oplus \Delta_{p}^{\prime \prime}$. The natural bundle isomorphism identifies $\Delta_{p}^{\prime \prime}=\Delta_{p-1}^{\prime}$. Since $a_{k}\left(x, T, \Delta_{p}^{M}\right)=a_{k}\left(x, T, \Delta_{p}^{\prime}\right)+a_{k}\left(x, T, \Delta_{p}^{\prime \prime}\right)=$ $a_{k}\left(x, T, \Delta_{p}^{\prime}\right)+a_{k}\left(x, T, \Delta_{p-1}^{\prime}\right)$, the alternating sum defining $a_{k}(x, T)_{\text {de Rham }}$ yields zero in this example which completes the proof.

We normalize the coordinate system $Z$ by requiring that $g_{i j}\left(x_{0}\right)=\delta_{i j}$. There are further normalizations we shall discuss shortly. There is a natural restriction map $r: \mathcal{T}_{m, n, k} \rightarrow \mathcal{T}_{m-1, n-1, k}$ defined algebraically by:

$$
\begin{aligned}
& r\left(g_{i j / \alpha}\right)= \begin{cases}0 & \text { if } \operatorname{deg}_{1}\left(g_{i j / \alpha}\right) \neq 0 \\
g_{i j / \alpha}^{\prime} & \text { if } \operatorname{deg}_{1}\left(g_{i j / \alpha}\right)=0\end{cases} \\
& r\left(T_{i / \beta}\right)= \begin{cases}0 & \text { if } \operatorname{deg}_{1}\left(T_{i / \beta}\right) \neq 0 \\
T_{i / \beta}^{\prime} & \text { if } \operatorname{deg}_{1}\left(T_{i / \beta}\right)=0\end{cases}
\end{aligned}
$$

In this expression, we let $g^{\prime}$ and $T^{\prime}$ denote the renumbered indices to refer to a manifold of one lower dimension. (It is inconvenient to have used the last $m-n$ indices for the normal bundle at this point, but again this is the cannonical convention which we have chosen not to change). In particular we note that Lemma 4.5.5 implies $r\left(a_{k}(x, T)_{\text {de Rham }}\right)=0$. Theorem 2.4.7 generalizes to this setting as:

Lemma 4.5.6. Let $p \in \mathcal{T}_{m, n, k}$ be such that $r(p)=0$.
(a) If $k$ is odd or if $k<n$ then $p=0$.
(b) If $k=n$ is even we let $E_{n}$ be the Euler form of the metric on $N$. Then $a_{n}(x, T)_{\text {de Rham }}=\left|\operatorname{det}\left(I-d T_{\nu}\right)\right|^{-1} f\left(d T_{\nu}\right) E_{n}$ for some $\mathrm{GL}(m-n)$ invariant smooth function $f(*)$.

We postpone the proof of this lemma for the moment to complete our discussion of the Lefschetz fixed point formula for the de Rham complex.

Theorem 4.5.7. Let $T: M \rightarrow M$ be non-degenerate with fixed point set consisting of the disjoint union of the submanifolds $N_{1}, N_{2}, \ldots, N_{r}$. Let $N$ denote one component of the fixed point set of dimension $n$ and let $a_{k}(x, T)_{\text {de Rham }}$ be the invariant of the heat equation. Then:
(a) $a_{k}(x, T)_{\text {de Rham }}=0$ for $k<n$ or if $k$ is odd.
(b) If $k=n$ is even, then $a_{n}(x, T)_{\text {de Rham }}=\operatorname{sign}\left(\operatorname{det}\left(I-d T_{\nu}\right)\right) E_{n}$.
(c) $L(T)_{\text {de Rham }}=\sum_{\nu} \operatorname{sign}(\operatorname{det}(I-d T)) \chi\left(N_{\nu}\right)$ (Classical Lefschetz fixed point formula).
Proof: (a) follows directly from Lemmas 4.5.4-4.5.7. We also conclude that if $k=n$ is even, then $a_{n}(x, T)_{\text {de Rham }}=h\left(d T_{\nu}\right) E_{n}$ for some invariant
function $h$. We know $h\left(d T_{\nu}\right)=\operatorname{sign}(\operatorname{det}(I-d T))$ if $n=m$ by Theorem 1.8.4. The multiplicative nature of the de Rham complex with respect to products establishes this formula in general which proves (b). Since

$$
L(T)_{\text {de Rham }}=\sum_{\mu} \operatorname{sign}\left(\operatorname{det}\left(I-d T_{\nu}\right)\right) \int_{N_{\mu}} E_{n_{\mu}} \operatorname{dvol}\left(x_{\mu}\right)
$$

by Theorem 1.8.2, (c) follows from the Gauss-Bonnet theorem already established. This completes the proof of Theorem 4.5.7.

Before beginning the proof of Lemma 4.5.6, we must further normalize the coordinates being considered. We have assumed that the metric chosen makes $N$ a totally geodesic submanifold. (In fact, this assumption is inessential, and the theorem is true in greater generality. The proof, however, is more complicated in this case). This implies that we can normalize the coordinate system chosen so that $g_{i j / k}\left(x_{0}\right)=0$ for $1 \leq i, j, k \leq m$. (If the submanifold is not totally geodesic, then the second fundamental form enters in exactly the same manner as it did for the Gauss-Bonnet theorem for manifolds with boundary).

By hypothesis, we decomposed $\left.T(M)\right|_{N}=T(N) \oplus \nu$ and decomposed $d T=I \oplus d T_{\nu}$. This implies that the Jacobian matrix $T_{i / j}$ satisfies:

$$
T_{a / b}=\delta_{a / b}, \quad T_{a / u}=T_{u / a}=0 \text { along } N
$$

Consequently:

$$
T_{a / \beta_{N} \beta_{\nu}}=0 \text { for }\left|\beta_{\nu}\right|<2 \quad \text { and } \quad T_{u / \beta_{N} \beta_{\nu}}=0 \text { for }\left|\beta_{\nu}\right|=0 .
$$

Consequently, for the non-zero $T_{i / \beta}$ variables:

$$
\sum_{a \leq n} \operatorname{deg}_{a}\left(T_{i / \beta}\right) \leq \operatorname{ord}\left(T_{i / \beta}\right)=|\beta|-1
$$

Let $p$ satisfy the hypothesis of Lemma 4.5.6. Let $p \neq 0$ and let $A$ be a monomial of $p$. We decompose:
$A=f\left(d T_{\nu}\right) \cdot A_{1} \cdot A_{2} \quad$ where $\quad \begin{cases}A_{1}=g_{i_{1} j_{1} / \alpha_{1}} \ldots g_{i_{r} j_{r} / \alpha_{r}} & \text { for }\left|\alpha_{\nu}\right| \geq 2 \\ A_{2}=A_{k_{1} / \beta_{1}} \ldots A_{k_{s} / \beta_{s}} & \text { for }\left|\beta_{\nu}\right| \geq 2 .\end{cases}$
Since $r(p)=0, \operatorname{deg}_{1} A \neq 0$. Since $p$ is invariant under orientation reversing changes of coordinates on $N, \operatorname{deg}_{1} A$ must be even. Since $p$ is invariant under coordinate permutations on $N, \operatorname{deg}_{a} A \geq 2$ for $1 \leq a \leq n$. We now count indices:

$$
\begin{aligned}
2 n & \leq \sum_{1 \leq a \leq n} \operatorname{deg}_{a}(A)=\sum_{1 \leq a \leq n} \operatorname{deg}_{a}\left(A_{1}\right)+\sum_{1 \leq a \leq n} \operatorname{deg}_{a}\left(A_{2}\right) \\
& \leq 2 r+\operatorname{ord}\left(A_{1}\right)+\operatorname{ord}\left(A_{2}\right) \leq 2 \operatorname{ord}\left(A_{1}\right)+2 \operatorname{ord}\left(A_{2}\right)=2 \operatorname{ord}(A)
\end{aligned}
$$

This shows $p=0$ if $\operatorname{ord}(A)<n$.
We now study the limiting case $\operatorname{ord}(A)=n$. All these inequalities must have been equalities. This implies in particular that ord $\left(A_{2}\right)=2 \operatorname{ord}\left(A_{2}\right)$ so $\operatorname{ord}\left(A_{2}\right)=0$ and $A_{2}$ does not appear. Furthermore, $4 r=\operatorname{ord}\left(A_{1}\right)$ implies $A_{1}$ is a polynomial in the 2-jets of the metric. Finally, $\sum_{1 \leq a \leq n} \operatorname{deg}_{a}\left(A_{1}\right)=$ $4 r$ implies that $A_{1}$ only depends upon the $\left\{g_{a b / c d}\right\}$ variab̄les. Since $p \neq 0$ implies $n=2 r$, we conclude $p=0$ if $k=n$ is odd. This completes the proof of (a).

To prove (b), we know that $p$ is a polynomial in the $\left\{g_{a b / c d}\right\}$ variables with coefficients which depend on the $\left\{T_{u / v}\right\}=d T_{\nu}$ variables. Exactly the same arguments as used in section 2.5 to prove Theorem 2.4.7 now show that $p$ has the desired form. This completes the proof; the normal and tangential indices decouple completely.

The Lefschetz fixed point formulas have been generalized by (Donnelly, The eta invariant of $G$-spaces) to the case of manifolds with boundary. We briefly summarize his work. Let $N$ be a compact Riemannian manifold with boundary $M$. We assume the metric is product near the boundary so a neighborhood of $M$ in $N$ has the form $M \times[0, \varepsilon)$. Let $n$ be the normal parameter. Let $G$ be a finite group acting on $N$ by isometries; $G$ must preserve $M$ as a set. We suppose $G$ has no fixed points on $M$ and only isolated fixed points in the interior of $N$ for $g \neq I$. Let $F(g)$ denote the set of fixed points for $g \neq I$. Let $\bar{M}=M / G$ be the resulting quotient manifold. It bounds a $V$-manifold in the sense of Kawasaki, but does not necessarily bound a smooth manifold as $N / G$ need not be a manifold.

Let $P: C^{\infty}\left(V_{1}\right) \rightarrow C^{\infty}\left(V_{2}\right)$ be an elliptic first order differential complex over $N$. Near the boundary, we assume $P$ has the form $P=p(d n)(\partial / \partial n+$ $A$ ) where $A$ is a self-adjoint elliptic first order differential operator over $M$ whose coefficients are independent of the normal parameter. We also assume the $G$ action on $N$ extends to an action on this elliptic complex. Then $g A=A g$ for all $g \in G$. Decompose $L^{2}\left(\left.V_{1}\right|_{M}\right)=\bigoplus_{\lambda} E(\lambda)$ into the finite dimensional eigenspaces of $A$. Then $g$ induces a representation on each $E(\lambda)$ and we define:

$$
\eta(s, A, g)=\sum_{\lambda} \operatorname{sign}(\lambda) \cdot|\lambda|^{-s} \operatorname{Tr}(g \text { on } E(\lambda))
$$

as the equivariant version of the eta invariant. This series converges absolutely for $\operatorname{Re}(s) \gg 0$ and has a meromorphic extension to $\mathbf{C}$. It is easy to see using the methods previously developed that this extension is regular for all values of $s$ since $g$ has no fixed points on $M$ for $g \neq I$. If $g=I$, this is just the eta invariant previously defined.

Let $B$ be orthonormal projection on the non-negative spectrum of $A$. This defines a non-local elliptic boundary value problem for the operator
$P$. Since $g$ commutes with the operator $A$, it commutes with the boundary conditions. Let $L(P, B, g)$ be the Lefschetz number of this problem. Theorems 4.3.11 and 1.8.3 generalize to this setting to become:

Theorem 4.5.8 (Donnelly). Let $N$ be a compact Riemannian manifold with boundary $M$. Let the metric on $N$ be product near $M$. Let $P: C^{\infty}\left(V_{1}\right) \rightarrow C^{\infty}\left(V_{2}\right)$ be a first order elliptic differential complex over $N$. Assume near $M$ that $P$ has the form $P=p(d n)(\partial / \partial n+A)$ where $A$ is a self-adjoint elliptic tangential operator on $C^{\infty}\left(V_{1}\right)$ over the boundary $M$. Let $L^{2}\left(\left.V_{1}\right|_{M}\right)=\bigoplus_{\lambda} E(\lambda)$ be a spectral resolution of $A$ and let $B$ be orthonormal projection on the non-negative spectrum of $A . \quad(P, B)$ is an elliptic boundary value problem. Assume given an isometry $g: N \rightarrow N$ with isolated fixed points $x_{1}, \ldots, x_{r}$ in the interior of $N$. Assume given an action $g$ on $V_{i}$ so $g P=P g$. Then $g A=A g$ as well. Define:

$$
\eta(s, A, g)=\sum_{\lambda} \operatorname{sign}(\lambda)|\lambda|^{-s} \cdot \operatorname{Tr}(g \text { on } E(\lambda))
$$

This converges for $\operatorname{Re}(s) \gg 0$ and extends to an entire function of $s$. Let $L(P, B, g)$ denote the Lefschetz number of $g$ on this elliptic complex and let

$$
\operatorname{defect}(P, g)\left(x_{i}\right)=\left\{\left(\operatorname{Tr}\left(g \text { on } V_{1}\right)-\operatorname{Tr}\left(g \text { on } V_{2}\right)\right) / \operatorname{det}(I-d T)\right\}\left(x_{i}\right)
$$

then:

$$
L(P, B, g)=\left\{\sum_{i} \operatorname{defect}(P, G)\left(x_{i}\right)\right\}-\frac{1}{2}\{\eta(0, A, g)+\operatorname{Tr}(g \text { on } \mathrm{N}(A))\}
$$

Remark: Donnelly's theorem holds in greater generality as one does not need to assume the fixed points are isolated and we refer to (Donnelly, The eta invariant of $G$-spaces) for details.

We use this theorem to compute the eta invariant on the quotient manifold $\bar{M}=M / G$. Equivariant eigensections for $A$ over $M$ correspond to the eigensections of $\bar{A}$ over $\bar{M}$. Then:

$$
\begin{aligned}
\tilde{\eta}(\bar{A})= & \frac{1}{2}\{\eta(0, \bar{A})+\operatorname{dim} \mathrm{N}(\bar{A})\} \\
= & \frac{1}{|G|} \cdot \frac{1}{2} \cdot \sum_{g \in G}\{\eta(0, A, g)+\operatorname{Tr}(g \text { on } \mathrm{N}(A))\} \\
= & \frac{1}{|G|} \sum_{g \in G}\{-L(P, B, g)\}+\frac{1}{|G|} \int_{N}\left(a_{n}\left(x, P^{*} P\right)-a_{n}\left(x, P P^{*}\right)\right) d x \\
& +\frac{1}{|G|} \sum_{\substack{g \in G \\
g \neq I}} \sum_{x \in F(g)} \operatorname{defect}(P, g)(x)
\end{aligned}
$$

The first sum over the group gives the equivariant index of $P$ with the given boundary condition. This is an integer and vanishes in $\mathbf{R} \bmod \mathbf{Z}$. The second contribution arises from Theorem 4.3.11 for $g=I$. The final contribution arises from Theorem 4.5.8 for $g \neq I$ and we sum over the fixed points of $g$ (which may be different for different group elements).

If we suppose $\operatorname{dim} M$ is even, then $\operatorname{dim} N$ is odd so $a_{n}\left(x, P^{*} P\right)-$ $a_{n}\left(x, P P^{*}\right)=0$. This gives a formula for $\eta(\bar{A})$ over $\bar{M}$ in terms of the fixed point data on $N$. By replacing $V_{i}$ by $V_{i} \otimes 1^{k}$ and letting $G$ act by a representation $\rho$ of $G$ in $\mathrm{U}(k)$, we obtain a formula for $\eta\left(\bar{A}_{\rho}\right)$.

If $\operatorname{dim} M$ is odd and $\operatorname{dim} N$ is even, the local interior formula for index $(P, B)$ given by the heat equation need not vanish identically. If we twist by a virtual representation $\rho$, we alter the defect formulas by multiplying by $\operatorname{Tr}(\rho(g))$. The contribution from Theorem 4.3.11 is multiplied by $\operatorname{Tr}(\rho(1))=\operatorname{dim} \rho$. Consequently this term disappears if $\rho$ is a representation of virtual dimension 0 . This proves:

Theorem 4.5.9. Let $N$ be a compact Riemannian manifold with boundary $M$. Let the metric on $N$ be product near $M$. Let $P: C^{\infty}\left(V_{1}\right) \rightarrow C^{\infty}\left(V_{2}\right)$ be a first order elliptic differential complex over $N$. Assume near $M$ that $P$ has the form $P=p(d n)(\partial / \partial+A)$ where $A$ is a self-adjoint elliptic tangential operator on $C^{\infty}\left(V_{1}\right)$ over $M$. Let $G$ be a finite group acting by isometries on $N$. Assume for $g \neq I$ that $g$ has only isolated fixed points in the interior of $N$, and let $F(g)$ denote the fixed point set. Assume given an action on $V_{i}$ so $g P=P g$ and $g A=A g$. Let $\bar{M}=M / G$ be the quotient manifold and $\bar{A}$ the induced self-adjoint elliptic operator on $C^{\infty}\left(\bar{V}_{1}=V_{1} / G\right)$ over $\bar{M}$. Let $\rho \in R(G)$ be a virtual representation. If $\operatorname{dim} M$ is odd, we assume $\operatorname{dim} \rho=0$. Then:

$$
\tilde{\eta}\left(\bar{A}_{\rho}\right)=\frac{1}{|G|} \sum_{\substack{g \in G \\ g \neq I}} \sum_{s \in F(g)} \operatorname{Tr}(\rho(g)) \operatorname{defect}(P, g)(x) \quad \bmod \mathbf{Z}
$$

for

$$
\operatorname{defect}(P, g)(x)=\left\{\left(\operatorname{Tr}\left(g \text { on } V_{1}\right)-\operatorname{Tr}\left(g \text { on } V_{2}\right)\right) / \operatorname{det}(I-d g)\right\}(x)
$$

We shall use this result in the next section and also in section 4.9 to discuss the eta invariant for sherical space forms.

### 4.6. The Eta Invariant and the $K$-Theory of Spherical Space Forms.

So far, we have used $K$-theory as a tool to prove theorems in analysis. $K$-theory and the Chern isomorphism have played an important role in our discussion of both the index and the twisted index theorems as well as in the regularity of eta at the origin. We now reverse the process and will use analysis as a tool to compute $K$-theory. We shall discuss the $K$-theory of spherical space forms using the eta invariant to detect the relevant torsion.

In section 4.5, Corollary 4.5.9, we discussed the equivariant computation of the eta invariant. We apply this to spherical space forms as follows. Let $G$ be a finite group and let $\tau: G \rightarrow \mathrm{U}(l)$ be a fixed point free representation. We suppose $\operatorname{det}(I-\tau(g)) \neq 0$ for $g \neq I$. Such a representation is necessarily faithful; the existence of such a representation places severe restrictions on the group $G$. In particular, all the Sylow subgroups for odd primes must be cyclic and the Sylow subgroup for the prime 2 is either cyclic or generalized quaternionic. These groups have all been classified by (Wolf, Spaces of Constant Curvature) and we refer to this work for further details on the subject.
$\tau(G)$ acts without fixed points on the unit sphere $S^{2 l-1}$ in $\mathbf{C}^{l}$. Let $M=$ $M(\tau)=S^{2 l-1} / \tau(G)$. We suppose $l>1$ so, since $S^{2 l-1}$ is simply connected, $\tau$ induces an isomorphism between $G$ and $\pi_{1}(M) . \quad M$ inherits a natural orientation and Riemannian metric. It also inherits a natural CauchyRiemann and $\mathrm{SPIN}_{c}$ structure. ( $M$ is not necessarily a spin manifold). The metric has constant positive sectional curvature. Such a manifold is called a spherical space form; all odd dimensional compact manifolds without boundary admitting metrics of constant positive sectional curvature arise in this way. The only even dimensional spherical space forms are the sphere $S^{2 l}$ and the projective space $\mathbf{R} P^{2 l}$. We concentrate for the moment on the odd dimensional case; we will return to consider $\mathbf{R} P^{2 l}$ later in this section.

We have the geometrical argument:

$$
T\left(S^{2 l-1}\right) \oplus 1=\left.T\left(\mathrm{R}^{2 l}\right)\right|_{S^{2 l-1}}=S^{2 l-1} \times \mathrm{R}^{2 l}=S^{2 l-1} \times \mathbf{C}^{l}
$$

is the trivial complex bundle of dimension $l$. The defining representation $\tau$ is unitary and acts naturally on this bundle. If $V_{\tau}$ is the locally flat complex bundle over $M(\tau)$ defined by the representation of $\pi_{1}(M(\tau))=G$, then this argument shows

$$
\left(V_{\tau}\right)_{\text {real }}=T(M(\tau)) \oplus 1 ;
$$

this is, of course, the Cauchy-Riemann structure refered to previously. In particular $V_{\tau}$ admits a nowhere vanishing section so we can split $V_{\tau}=V_{1} \oplus 1$ where $V_{1}$ is an orthogonal complement of the trivial bundle corresponding to the invariant normal section of $\left.T\left(\mathrm{R}^{2 l}\right)\right|_{S^{2 l-1}}$. Therefore

$$
\sum_{\nu}(-1)^{\nu} \Lambda^{\nu}\left(V_{\tau}\right)=\sum_{\nu}(-1)^{\nu} \Lambda^{\nu}\left(V_{1} \oplus 1\right)=\sum_{\nu}(-1)^{\nu}\left\{\Lambda^{\nu}\left(V_{1}\right) \oplus \Lambda^{\nu-1}\left(V_{1}\right)\right\}=0
$$

in $K(M)$. This bundle corresponds to the virtual representation $\alpha=$ $\sum_{\nu}(-1)^{\nu} \Lambda^{\nu}(\tau) \in R_{0}(G)$. This proves:
Lemma 4.6.1. Let $\tau$ be a fixed point free representation of a finite group $G$ in $\mathrm{U}(l)$ and let $M(\tau)=S^{2 l-1} / \tau(G)$. Let $\alpha=\sum_{\nu}(-1)^{\nu} \Lambda^{\nu}(\tau) \in R_{0}(G)$. Then:

$$
T(M(\tau)) \oplus 1=\left(V_{\tau}\right)_{\text {real }} \quad \text { and } \quad V_{\alpha}=0 \text { in } \widetilde{K}(M(\tau))
$$

The sphere bounds a disk $D^{2 l}$ in $C^{l}$. The metric is not product near the boundary of course. By making a radial change of metric, we can put a metric on $D^{2 l}$ agreeing with the standard metric at the origin and with a product metric near the boundary so that the action of $\mathrm{O}(2 l)$ continues to be an action by isometries. The transition functions of $V_{\tau}$ are unitary so $T(M(\tau))$ inherits a natural $\mathrm{SPIN}_{c}$ structure. Let $*=$ signature or Dolbeault and let $A_{*}$ be the tangential operator of the appropriate elliptic complex over the disk. $\tau(G)$ acts on $D^{2 l}$ and the action extends to an action on both the signature and Dolbeault complexes. There is a single fixed point at the origin of the disk. Let $\operatorname{defect}(\tau(g), *)$ denote the appropriate term from the Lefschetz fixed point formulas. Let $\left\{\lambda_{\nu}\right\}$ denote the complex eigenvalues of $\tau(g)$, and let $\tau(g)_{r}$ denote the corresponding element of $\mathrm{SO}(2 l)$. It follows from section 4.5 that:

$$
\begin{aligned}
& \operatorname{defect}(\tau(g), \text { signature })=\prod_{\nu} \frac{\lambda_{\nu}+1}{\lambda_{\nu}-1} \\
& \operatorname{defect}(\tau(g), \text { Dolbeault })=\frac{\operatorname{det}(\tau(g))}{\operatorname{det}(\tau(g)-I)}=\prod_{\nu} \frac{\lambda_{\nu}}{\lambda_{\nu}-1}
\end{aligned}
$$

We apply Corollary 4.5.9 to this situation to compute:
4.6.2. Let $\tau: G \rightarrow \mathrm{U}(l)$ be a fixed point free representation of a finite group. Let $M(\tau)=S^{2 l-1} / \tau(G)$ be a spherical space form. Let $*=$ signature or Dolbeault and let $A_{*}$ be the tangential operator of the appropriate elliptic complex. Let $\rho \in R_{0}(G)$ be a virtual representation of dimension 0. Then:

$$
\tilde{\eta}\left(\left(A_{*}\right)_{\rho}\right)=\frac{1}{|G|} \sum_{\substack{g \in G \\ g \neq I}} \operatorname{Tr}(\rho(g)) \operatorname{defect}(\tau(g), *) \quad \text { in } \mathbf{R} \bmod \mathbf{Z}
$$

Remark: A priori, this identity is in $\mathbf{R} \bmod \mathbf{Z}$. It is not difficult to show that this generalized Dedekind sum is always $\mathbf{Q} \bmod \mathbf{Z}$ valued and that $|G|^{l} \tilde{\eta} \in \mathbf{Z}$ so one has good control on the denominators involved.

The perhaps somewhat surprising fact is that this invariant is polynomial. Suppose $G=\mathbf{Z}_{n}$ is cyclic. Let $x=\left(x_{1}, \ldots, x_{l}\right)$ be a collection of
indeterminates. Let $T d$ and $L$ be the Todd and Hirzebruch polynomials discussed previously. We define:

$$
\begin{array}{rlrl}
T d_{0}(\vec{x}) & =1 & L_{0}(\vec{x}) & =1 \\
T d_{1}(\vec{x}) & =\frac{1}{2} \sum x_{j} & L_{1}(\vec{x}) & =0 \\
T d_{2}(\vec{x}) & =\frac{1}{12}\left\{\sum_{j<k} x_{j} x_{k}+\left(\sum_{j} x_{j}\right)^{2}\right\} & L_{2}(\vec{x})=\frac{1}{3} \sum x_{i}^{2}
\end{array}
$$

where we have renumbered the $L_{k}$ polynomials to be homogeneous of degree $k$. Let $s$ be another parameter which represents the first Chern class of a line bundle. The integrands of the index formula are given by:

$$
\begin{aligned}
P_{l}(s ; \vec{x} ; \text { signature }) & =\sum_{j+k=l} s^{j} L_{k}(\vec{x}) 2^{j} / j! \\
P_{l}(s ; \vec{x} ; \text { Dolbeault }) & =\sum_{j+k=l} s^{j} T d_{k}(\vec{x}) / j!
\end{aligned}
$$

Let $\mu(l)$ denote the least common denominator of these rational polynomials.

We identify $\mathbf{Z}_{n}$ with the group of $n^{\text {th }}$ roots of unity in C. Let $\rho_{s}(\lambda)=$ $\lambda^{s}$ for $0 \leq s<n$ parameterize the irreducible representations. If $\vec{q}=$ $\left(q_{1}, \ldots, q_{l}\right)$ is a collection of integers coprime to $n$, let $\tau=\rho_{q_{1}} \oplus \cdots \oplus \rho_{q_{l}}$ so $\tau(\lambda)=\operatorname{diag}\left(\lambda^{q_{1}}, \ldots, \lambda^{q_{l}}\right)$. This is a fixed pont free representation; up to unitary equivalence any fixed point free representation has this form. Let $L(n ; \vec{q})=M(\tau)=S^{2 l-1} / \tau\left(\mathbf{Z}_{n}\right)$ be the corresponding spherical space form; this is called a lens space.
Lemma 4.6.3. Let $M=L(n ; \vec{q})$ be a lens space of dimension $2 l-1$. Let $e \in \mathbf{Z}$ satisfy $e q_{1} \ldots q_{l} \equiv 1 \bmod n \cdot \mu(l)$. Let $*=$ signature or Dolbeault and let $A_{*}$ be the tangential operator of the corresponding elliptic complex. Let $P_{l}(s ; x ; *)$ denote the corresponding rational polynomial as defined above. Then

$$
\operatorname{ind}\left(\rho_{s}-\rho_{0}, A_{*}\right)=-\frac{e}{n}\left\{P_{l}(s ; n, \vec{q} ; *)-P_{l}(0 ; n, \vec{q} ; *)\right\} \bmod \mathbf{Z}
$$

Remark: If $M$ admits a spin structure, there is a corresponding formula for the tangential operator of the spin complex. This illustrates the close relationship between the Lefschetz fixed point formulas, the Atiyah-Singer index theorem, and the eta invariant, as it ties together all these elements. If $l=2$ so $M=L(n ; 1, q)$ then

$$
\operatorname{ind}\left(\rho_{s}-\rho_{0}, A_{\text {signature }}\right) \equiv-\frac{q^{\prime}}{n} \cdot(2 s)^{2} \cdot \frac{1}{2} \equiv \frac{-q^{\prime} \cdot 2 s^{2}}{n}
$$

which is the formula obatined previously in section 4.4.
Proof: We refer to (Gilkey, The eta invariant and the $K$-theory of odd dimensional spherical space forms) for the proof; it is a simple residue calculation using the results of Hirzebruch-Zagier.

This is a very computable invariant and can be calculated using a computer from either Lemma 4.6.2 or Lemma 4.6.3. Although Lemma 4.6.3 at first sight only applies to cyclic groups, it is not difficult to use the Brauer induction formula and some elementary results concerning these groups to obtain similar formulas for arbitrary finite groups admitting fixed point free representations.

Let $M$ be a compact manifold without boundary with fundamental group $G$. If $\rho$ is a unitary representation of $G$ and if $P$ is a self-adjoint elliptic differential operator, we have defined the invariant $\tilde{\eta}\left(P_{\rho}\right)$ as an $\mathbf{R} \bmod \mathbf{Z}$ valued invariant. (In fact this invariant can be defined for arbitrary representations and for elliptic pseudo-differential operators with leading symbol having no purely imaginary eigenvalues on $S\left(T^{*} M\right)$ and most of what we will say will go over to this more general case. As we are only interested in finite groups it suffices to work in this more restricted category).

Let $R(G)$ be the group representation ring of unitary virtual representations of $G$ and let $R_{0}(G)$ be the ideal of virtual representations of virtual dimension 0 . The map $\rho \mapsto V_{\rho}$ defines a ring homomorphism from $R(G)$ to $K(M)$ and $R_{0}(G)$ to $\widetilde{K}(M)$. We shall denote the images by $K_{\text {flat }}(M)$ and $\widetilde{K}_{\text {flat }}(M)$; these are the rings generated by virtual bundles admitting locally flat structures, or equivalently by virtual bundles with constant transition functions. Let $P$ be a self-adjoint and elliptic differential oprator. The map $\rho \mapsto \tilde{\eta}\left(A_{\rho}\right)$ is additive with respect to direct sums and extends to a $\operatorname{map} R(G) \rightarrow \mathbf{R} \bmod \mathbf{Z}$ as already noted. We let $\operatorname{ind}(\rho, P)$ be the map form $R_{0}(G)$ to $\mathbf{R} \bmod \mathbf{Z}$. This involves a slight change of notation from section 4.4; if $\rho$ is a representation of $G$, then

$$
\operatorname{ind}(\rho-\operatorname{dim}(\rho) \cdot 1, P)
$$

denotes the invariant previously defined by $\operatorname{ind}(\rho, P)$. This invariant is constant under deformations of $P$ within this class; the Atiyah-Patodi-Singer index theorem for manifolds with boundary (Theorem 4.3.11) implies it is also an equivariant cobordism invariant. We summarize its relevant properties:

Lemma 4.6.4. Let $M$ be a compact smooth manifold without boundary.
Let $P$ be a self-adjoint elliptic differential operator over $M$ and let $\rho \in$ $R_{0}\left(\pi_{1}(M)\right)$.
(a) Let $P(a)$ be a smooth 1-parameter family of such operators, then $\operatorname{ind}(\rho, P(a))$ is independent of $a$ in $\mathbf{R} \bmod \mathbf{Z}$. If $G$ is finite, this is $\mathbf{Q} \bmod \mathbf{Z}$ valued.
(b) Let $P$ be a first order operator. Suppose there exists a compact manifold $N$ with $d N=M$. Suppose there is an elliptic complex $Q: C^{\infty}\left(V_{1}\right) \rightarrow$ $C^{\infty}\left(V_{2}\right)$ over $N$ so $P$ is the tangential part of $Q$. Suppose the virtual bundle $V_{\rho}$ can be extended as a locally flat bundle over $N$. Then $\operatorname{ind}(\rho, P)=0$.

Proof: The first assertion of (a) follows from Lemma 4.4.1. The locally flat bundle $V_{\rho}$ is rationally trivial so by multipying by a suitable integer we can actually assume $V_{\rho}$ corresponds to a flat structure on the difference of trivial bundles. The index is therefore given by a local formula. If we lift to the universal cover, we multiply this local formula by $|G|$. On the universal cover, the index vanishes identically as $\pi_{1}=0$. Thus an integer multiple of the index is 0 in $\mathbf{R} / \mathbf{Z}$ so the index is in $\mathbf{Q} / \mathbf{Z}$ which proves (a). To prove (b) we take the operator $Q$ with coefficients in $V_{\rho}$. The local formula of the heat equation is just multiplied by the scaling constant $\operatorname{dim}(\rho)=0$ since $V_{\rho}$ is locally flat over $N$. Therefore Theorem 4.3.11 yields the identity index $\left(Q, B\right.$, coeff in $\left.V_{\rho}\right)=0-\operatorname{ind}(\rho, P)$. As the index is always an integer, this proves (b). We will use (b) in section 4.9 to discuss isospectral manifolds which are not diffeomorphic.

Examle 4.4 .2 shows the index is not an invariant in $K$-theory. We get $K$-theory invariants as follows:
Lemma 4.6.5. Let $M$ be a compact manifold without boundary. Let $P$ be an elliptic self-adjoint pseudo-differential operator. Let $\rho_{\nu} \in R_{0}\left(\pi_{1}(M)\right)$ and define the associative symmetric bilinear form on $R_{0} \otimes R_{0}$ by:

$$
\operatorname{ind}\left(\rho_{1}, \rho_{2}, P\right)=\operatorname{ind}\left(\rho_{1} \otimes \rho_{2}, P\right)
$$

This takes values in $\mathbf{Q} \bmod \mathbf{Z}$ and extends to an associative symmetric bilinear form $\operatorname{ind}(*, *, P): \widetilde{K}_{\text {flat }}(M) \otimes \widetilde{K}_{\text {flat }}(M) \rightarrow \mathbf{Q} \bmod \mathbf{Z}$. If we consider the dependence upon $P$, then we get a trilinear map

$$
\text { ind: } \widetilde{K}_{\text {flat }}(M) \otimes \widetilde{K}_{\text {flat }}(M) \otimes K\left(S\left(T^{*} M\right)\right) / K(M) \rightarrow \mathbf{Q} \bmod \mathbf{Z}
$$

Proof: The interpretation of the dependence in $P$ as a map in $K$-theory follows from 4.3.3 and is therefore omitted. Any virtual bundle admitting a locally flat structure has vanishing rational Chern character and must be a torsion class. Once we have proved the map extends to $K$-theory, it will follow it must be $\mathbf{Q} \bmod \mathbf{Z}$ valued. We suppose given representations $\rho_{1}, \hat{\rho}_{1}, \rho_{2}$ and a bundle isomorphism $V_{\rho_{1}}=V_{\hat{\rho}_{1}}$. Let $j=\operatorname{dim}\left(\rho_{1}\right)$ and $k=\operatorname{dim}\left(\rho_{2}\right)$. If we can show

$$
\operatorname{ind}\left(\left(\rho_{1}-j\right) \otimes\left(\rho_{2}-k\right), P\right)=\operatorname{ind}\left(\left(\hat{\rho}_{1}-j\right) \otimes\left(\rho_{2}-k\right), P\right)
$$

then the form will extend to $\widetilde{K}_{\text {flat }}(M) \otimes \widetilde{K}_{\text {flat }}(M)$ and the lemma will be proved.

We calculate that:

$$
\begin{aligned}
\operatorname{ind}\left(( \rho _ { 1 } - j ) \otimes \left(\rho_{2}\right.\right. & -k), P) \\
& =\tilde{\eta}\left(P_{\rho_{1} \otimes \rho_{2}}\right)+j \cdot k \cdot \tilde{\eta}(P)-j \cdot \tilde{\eta}\left(P_{\rho_{2}}\right)-k \cdot \tilde{\eta}\left(P_{\rho_{1}}\right) \\
& =\tilde{\eta}\left\{\left(P_{\rho_{1}}\right)_{\rho_{2}}\right\}-k \cdot \tilde{\eta}\left(P_{\rho_{1}}\right)-j \cdot\left\{\tilde{\eta}\left(P_{\rho_{2}}\right)-k \cdot \tilde{\eta}(P)\right\} \\
& =\operatorname{ind}\left(\rho_{2}, P_{\rho_{1}}\right)-j \cdot \operatorname{ind}\left(\rho_{2}, P\right)
\end{aligned}
$$

By hypothesis the bundles defined by $\rho_{1}$ and $\hat{\rho}_{1}$ are isomorphic. Thus the two operators $P_{\rho_{1}}$ and $P_{\hat{\rho}_{1}}$ are homotopic since they have the same leading symbol. Therefore $\operatorname{ind}\left(\rho_{2}, P_{\rho_{1}}\right)=\operatorname{ind}\left(\rho_{2}, P_{\hat{\rho}_{1}}\right)$ which completes the proof. We remark that this bilinear form is also associative with respect to multiplication by $R(G)$ and $K_{\text {flat }}(M)$.

We will use this lemma to study the $K$-theory of spherical space forms.
Lemma 4.6.6. Let $\tau$ be a fixed point free representation of a finite group G. Define

$$
\operatorname{ind}_{\tau}\left(\rho_{1}, \rho_{2}\right)=\frac{1}{|G|} \sum_{\substack{g \in G \\ g \neq I}} \operatorname{Tr}\left(\rho_{1} \otimes \rho_{2}\right)(g) \frac{\operatorname{det}(\tau(g))}{\operatorname{det}(\tau(g)-I)}
$$

Let $\alpha=\sum_{\nu}(-1)^{\nu} \Lambda^{\nu}(\tau) \in R_{0}(G)$. Let $\rho_{1} \in R_{0}(G)$ and suppose $\operatorname{ind}_{\tau}\left(\rho_{1}, \rho_{2}\right)$ $=0$ in $\mathbf{Q} \bmod \mathbf{Z}$ for all $\rho_{2} \in R_{0}(G)$. Then $\rho_{2} \in \alpha R(G)$.
Proof: The virtual representation $\alpha$ is given by the defining relation that
$\operatorname{Tr}(\alpha(g))=\operatorname{det}(I-\tau(g)) . \quad \operatorname{det}(\tau)$ defines a 1-dimensional representation of $G$; as this is an invertible element of $R(G)$ we see that the hypothesis implies

$$
\frac{1}{|G|} \sum_{\substack{g \in G \\ g \neq I}} \operatorname{Tr}\left(\rho_{1} \otimes \rho_{2}\right) / \operatorname{det}(I-\tau(g)) \in \mathbf{Z} \quad \text { for all } \rho_{2} \in R_{0}(G)
$$

If $f$ and $\tilde{f}$ are any two class functions on $G$, we define the symmetric inner product $(f, \tilde{f})=\frac{1}{|G|} \sum_{g \in G} f(g) \tilde{f}(g)$. The orthogonality relations show that $f$ is a virtual character if and only if $(f, \operatorname{Tr}(\rho)) \in \mathbf{Z}$ for all $\rho \in R(\mathbf{Z})$. We define:

$$
\begin{aligned}
& f(g)=\operatorname{Tr}\left(\rho_{1}(g)\right) / \operatorname{det}(I-\tau(g)) \quad \text { for } g \neq I \\
& f(g)=-\sum_{\substack{h \in G \\
h \neq I}} \operatorname{Tr}\left(\rho_{1}(h)\right) / \operatorname{det}(I-\tau(h)) \quad \text { if } g=I
\end{aligned}
$$

Then $(f, 1)=0$ by definition. As $\left(f, \rho_{2}\right) \in \mathbf{Z}$ by hypothesis for $\rho_{2} \in R_{0}(G)$ we see $\left(f, \rho_{2}\right) \in \mathbf{Z}$ for all $\rho_{2} \in R(G)$ so $f$ is a virtual character. We let $\operatorname{Tr}(\rho)(g)=f(g)$. The defining equation implies:

$$
\operatorname{Tr}(\rho \otimes \alpha)(g)=\operatorname{Tr}\left(\rho_{1}(g)\right) \quad \text { for all } g \in G
$$

This implies $\rho_{1}=\rho \otimes \alpha$ and completes the proof.
We can now compute the $K$-theory of odd dimensional spherical space forms.

THEOREM 4.6.7. Let $\tau: G \rightarrow \mathrm{U}(l)$ be a fixed point free representation of a finite group $G$. Let $M(\tau)=S^{2 l-1} / \tau(G)$ be a spherical space form. Suppose $l>1$. Let $\alpha=\sum_{\nu}(-1)^{\nu} \Lambda^{\nu}(\tau) \in R_{0}(G)$. Then $\widetilde{K}_{\text {flat }}(M)=R_{0}(G) / \alpha R(G)$ and

$$
\operatorname{ind}\left(*, *, A_{\text {Dolbeault }}\right): \widetilde{K}_{\text {flat }}(M) \otimes \widetilde{K}_{\text {flat }}(M) \rightarrow \mathbf{Q} \bmod \mathbf{Z}
$$

is a non-singular bilinear form.
Remark: It is a well known topological fact that for such spaces $\widetilde{K}=\widetilde{K}_{\text {flat }}$ so we are actually computing the reduced $K$-theory of odd dimensional spherical space forms (which was first computed by Atiyah). This particular proof gives much more information than just the isomorphism and we will draw some corollaries of the proof.

Proof: We have a surjective map $R_{0}(G) \rightarrow \widetilde{K}_{\text {flat }}(M)$. By Lemma 4.6.1 we have $\alpha \mapsto 0$ so $\alpha R(G)$ is in the kernel of the natural map. Conversely, suppose $V_{\rho}=0$ in $\widetilde{K}$. By Lemma 4.6.5 we have ind $\left(\rho, \rho_{1}, A_{\text {Dolbeault }}\right)=0$ for all $\rho_{1} \in R_{0}(G)$. Lemma 4.6.2 lets us identify this invariant with $\operatorname{ind}_{\tau}\left(\rho, \rho_{1}\right)$. Lemma 4.6.6 lets us conclude $\rho \in \alpha R(G)$. This shows the kernel of this map is precisely $\alpha R(G)$ which gives the desired isomorphism. Furthermore, $\rho \in \operatorname{ker}\left(\operatorname{ind}_{\tau}(\rho, *)\right)$ if and only if $\rho \in \alpha R(G)$ if and only if $V_{\rho}=0$ so the bilinear form is non-singular on $\widetilde{K}$.

It is possible to prove a number of other results about the $K$-theory ring using purely group theoretic methods; the existence of such a non-singular associative symmetric $\mathbf{Q} \bmod \mathbf{Z}$ form is an essential ingredient.

Corollary 4.6.8. Adopt the notation of Theorem 4.6.7.
(a) $\widetilde{K}(M)$ only depends on $(G, l)$ as a ring and not upon the particular $\tau$ chosen.
(b) The index of nilpotency for this ring is at most l-i.e., if $\rho_{\nu} \in R_{0}(G)$ then $\prod_{1 \leq \nu \leq l} \rho_{\nu} \in \alpha R_{0}(G)$ so the product of $l$ virtual bundles of $\widetilde{K}(M)$ always gives 0 in $\widetilde{K}(M)$.
(c) Let $V \in \widetilde{K}(M)$. Then $V=0$ if and only if $\pi^{*}(V)=0$ for all possible covering projections $\pi: L(n ; \vec{q}) \rightarrow M$ by lens spaces.

There is, of course, a great deal known concerning these rings and we refer to (N. Mahammed, $K$-theorie des formes spheriques) for further details.

If $G=\mathbf{Z}_{2}$, then the resulting space is $\mathbf{R} P^{2 l-1}$ which is projective space. There are two inequivalent unitary irreducible representations $\rho_{0}, \rho_{1}$ of $G$. Let $x=\rho_{1}-\rho_{0}$ generate $R_{0}\left(\mathbf{Z}_{2}\right)=\mathbf{Z}$; the ring structure is given by
$x^{2}=-2 x$. Let $A$ be the tangential operator of the Dolbeault complex:

$$
\begin{aligned}
\operatorname{ind}(x, A) & =\frac{1}{2} \operatorname{Tr}(x(-1)) \cdot \operatorname{det}\left(-I_{l}\right) / \operatorname{det}\left((-I)_{l}-I_{l}\right) \\
& =-2^{-l}
\end{aligned}
$$

using Lemma 4.6.2. Therefore $\operatorname{ind}(x, x, A)=2^{-l+1}$ which implies that $\widetilde{K}\left(\mathbf{R} P^{2 l-1}\right)=\mathbf{Z} / 2^{l-1} \mathbf{Z}$. Let $L=V_{\rho_{1}}$ be the tautological bundle over $\mathbf{R} P^{2 l-1}$. It is $S^{2 l-1} \times \mathbf{C}$ modulo the relation $(x, z)=(-x,-z)$. Using Clifford matrices we can construct a map $e: S^{2 l-1} \rightarrow \mathrm{U}\left(2^{l-1}\right)$ such that $e(-x)=-e(x)$. This gives an equivariant trivialization of $S^{2 l-1} \times \mathbf{C}^{2^{l-1}}$ which descends to give a trivialization of $2^{l-1} \cdot L$. This shows explicitly that $2^{l-1}(L-1)=0$ in $K\left(\mathbf{R} P^{2 l-1}\right)$; the eta invariant is used to show no lower power suffices. This proves:

Corollary 4.6.9. Let $M=\mathbf{R} P^{2 l-1}$ then $\widetilde{K}(M) \simeq \mathbf{Z} / 2^{l-1} \mathbf{Z}$ where the ring structure is $x^{2}=-2 x$.

Let $l=2$ and let $G=\mathbf{Z}_{n}$ be cyclic. Lemma 4.6.3 shows

$$
\operatorname{ind}\left(\rho_{s}-\rho_{0}, A_{\text {Dolbeault }}\right)=\frac{-q^{\prime}}{n} \cdot \frac{s^{2}}{2}
$$

Let $x=\rho_{1}-\rho_{0}$ so $x^{2}=\rho_{2}-\rho_{0}-2\left(\rho_{1}-\rho_{0}\right)$ and

$$
\operatorname{ind}\left(x, x, A_{\text {Dolbeault }}\right)=\frac{-q^{\prime}}{n} \cdot \frac{4-2}{2}
$$

is a generator of $\mathbf{Z}\left[\frac{1}{n}\right] \bmod \mathbf{Z}$. As $x^{2}=0$ in $\widetilde{K}$ and as $x \cdot R(G)=R_{0}(G)$ we see:

Corollary 4.6.10. Let $M=L(n ; 1, q)$ be a lens space of dimension 3. $\widetilde{K}(M)=\mathbf{Z}_{n}$ with trivial ring structure.

We have computed the $K$-theory for the odd dimensional spherical space forms. $\widetilde{K}\left(S^{2 l}\right)=\mathbf{Z}$ and we gave a generator in terms of Clifford algebras in Chapter 3. To complete the discussion, it suffices to consider even dimensional real projective space $M=\mathbf{R} P^{2 l}$. As $\widetilde{H}^{\text {even }}(M ; \mathbf{Q})=0, \widetilde{K}$ is pure torsion by the Chern isomorphism. Again, it is known that the flat and regular $K$-theory coincide. Let $x=L-1=V_{\rho_{1}}-V_{\rho_{0}}$. This is the restriction of an element of $\widetilde{K}\left(\mathbf{R} P^{2 l+1}\right)$ so $2^{l} x=0$ by Corollary 4.6.9. It is immediate that $x^{2}=-2 x$. We show $\widetilde{K}\left(\mathbf{R} P^{2 l+1}\right)=\mathbf{Z} / 2^{l} \mathbf{Z}$ by giving a surjective map to a group of order $2^{l}$.

We construct an elliptic complex $Q$ over the disk $D^{2 l+1}$. Let $\left\{e_{0}, \ldots, e_{2 l}\right\}$ be a collection of $2^{l} \times 2^{l}$ skew-adjoint matrices so $e_{j} e_{k}+e_{k} e_{j}=-2 \delta_{j k}$. Up to unitary equivalence, the only invariant of such a collection is $\operatorname{Tr}\left(e_{0} \ldots e_{2 l}\right)=$
$\pm(2 i)^{l}$. There are two inequivalent collections; the other is obtained by taking $\left\{-e_{0}, \ldots,-e_{2 l}\right\}$. Let $Q$ be the operator $\sum_{j} e_{j} \partial / \partial x_{j}$ acting to map $Q: C^{\infty}\left(V_{1}\right) \rightarrow C^{\infty}\left(V_{2}\right)$ where $V_{i}$ are trivial bundles of dimension $2^{l}$ over the disk. Let $g(x)=-x$ be the antipodal map. Let $g$ act by +1 on $V_{1}$ and by -1 on $V_{2}$, then $g Q=Q g$ so this is an equivariant action. Let $A$ be the tangential operator of this complex on $S^{2 l}$ and $\bar{A}$ the corresponding operator on $S^{2 l} / \mathbf{Z}_{2}=M . \bar{A}$ is a self-adjoint elliptic first order operator on $C^{\infty}\left(1^{2^{l}}\right)$. If we replace $Q$ by $-Q$, the tangential operator is unchanged so $\bar{A}$ is invariantly defined independent of the choice of the $\left\{e_{j}\right\}$. (In fact, $\bar{A}$ is the tangential operator of the $\mathrm{PIN}_{c}$ complex.)

The dimension is even so we can apply Corollary 4.5.9 to conclude

$$
\tilde{\eta}(\bar{A})=\frac{1}{2} \cdot \operatorname{det}\left(I_{2 l+1}-\left(-I_{2 l+1}\right)\right)^{-1}\left\{2^{l}-\left(-2^{l}\right)\right\}=2^{-l-1} .
$$

If we interchange the roles of $V_{1}$ and $V_{2}$, we change the sign of the eta invariant. This is equivalent to taking coefficients in the bundle $L=V_{\rho_{1}}$. Let $x=L-1$ then $\operatorname{ind}\left(\rho_{1}-\rho_{0}, \bar{A}\right)=-2^{-l-1}-2^{-l-1}=-2^{-l}$.

The even dimensional rational cohomology of $S\left(T^{*} M\right)$ is generated by $H^{0}$ and thus $K\left(S\left(T^{*} M\right)\right) / K(M) \otimes \mathbf{Q}=0$. Suppose $2^{l-1} \cdot x=0$ in $\widetilde{K}(M)$. Then there would exist a local formula for $2^{l-1} \operatorname{ind}\left(\rho_{1}-\rho_{0}, *\right)$ so we could lift this invariant from $\mathbf{Q} \bmod \mathbf{Z}$ to $\mathbf{Q}$. As this invariant is defined on the torsion group $K\left(S\left(T^{*} M\right)\right) / K(M)$ it would have to vanish. As $2^{l-1} \operatorname{ind}\left(\rho_{1}-\right.$ $\left.\rho_{0}, \bar{A}\right)=-\frac{1}{2}$ does not vanish, we conclude $2^{l-1} \cdot x$ is non-zero in $K$-theory as desired. This proves:
Corollary 4.6.11. $\widetilde{K}\left(\mathbf{R} P^{2 l}\right) \simeq \mathbf{Z} / 2^{l} \mathbf{Z}$. If $x=L-1$ is the generator, then $x^{2}=-2 x$.

We can squeeze a bit more out of this construction to compute the $K$-theory of the unit sphere bundle $K\left(S\left(T^{*} M\right)\right.$ ) where $M$ is a spherical space form. First suppose $\operatorname{dim} M=2 l-1$ is odd so $M=S^{2 l-1} / \tau(G)$ where $\tau$ is a unitary representation. Then $M$ has a Cauchy-Riemann structure and we can decompose $T^{*}(M)=1 \oplus V$ where $V$ admits a complex structure; $T(M) \oplus 1=\left(V_{\tau}\right)_{\text {real }}$. Thus $S\left(T^{*} M\right)$ has a non-vanishing section and the exact sequence $0 \rightarrow K(M) \rightarrow K\left(S\left(T^{*} M\right)\right) \rightarrow K\left(S\left(T^{*} M\right)\right) / K(M) \rightarrow 0$ splits. The usual clutching function construction permits us to identify $K\left(S\left(T^{*} M\right)\right) / K(M)$ with $K(V)$. As $V$ admits a complex structure, the Thom isomorphism identifies $K(V)=x \cdot K(M)$ where $x$ is the Thom class. This gives the structure $K\left(S\left(T^{*} M\right)\right)=K(M) \oplus x K(M)$. The bundle $x$ over $S\left(T^{*} M\right)$ can be taken to be $\Pi_{+}(p)$ where $p$ is the symbol of the tangential operator of the Dolbeault complex. The index form can be regarded as a pairing $\widetilde{K}(M) \otimes x \cdot \widetilde{K}(M) \rightarrow \mathbf{Q} \bmod \mathbf{Z}$ which is non-degenerate.

It is more difficult to analyse the even dimensional case. We wish to compute $K\left(S\left(T^{*} \mathbf{R} P^{2 l}\right)\right)$. $2^{l} \operatorname{ind}\left(\rho_{1}-\rho_{0}, *\right)$ is given by a local formula
since $2^{l} x=0$. Thus this invariant must be zero as $K\left(S\left(T^{*} M\right)\right) / K(M)$ is torsion. We have constructed an operator $\bar{A}$ so $\operatorname{ind}\left(\rho_{1}-\rho_{0}, \bar{A}\right)=-2^{-l}$ and thus $\operatorname{ind}\left(\rho_{1}-\rho_{0}, *\right): K\left(S\left(T^{*} M\right)\right) / K(M) \rightarrow \mathbf{Z}\left[2^{-l}\right] / \mathbf{Z}$ is surjective with the given range. This is a cyclic group of order $2^{l}$ so equivalently

$$
\operatorname{ind}\left(\rho_{1}-\rho_{0}, *\right): K\left(S\left(T^{*} M\right)\right) / K(M) \rightarrow \mathbf{Z} / 2^{l} \mathbf{Z} \rightarrow 0
$$

Victor Snaith (private communication) has shown us that the existence of such a sequence together with the Hirzebruch spectral sequence in $K$-theory shows

$$
0 \rightarrow K(M) \rightarrow K\left(S\left(T^{*} M\right)\right) \text { is exact and }\left|K\left(S\left(T^{*} M\right)\right) / K(M)\right|=2^{l}
$$

so that index $\left(\rho_{1}-\rho_{0}, *\right)$ becomes part of a short exact sequence:

$$
0 \rightarrow K(M) \rightarrow K\left(S\left(T^{*} M\right)\right) \rightarrow \mathbf{Z} / 2^{l} \mathbf{Z} \rightarrow 0
$$

To compute the structure of $K\left(S\left(T^{*} M\right)\right)=\mathbf{Z} \oplus \widetilde{K}\left(S\left(T^{*} M\right)\right)$, we must determine the group extension.

Let $x=L-1=V_{\rho_{1}}-V_{\rho_{0}}$ generate $\widetilde{K}(M)$. Let $\left.y=\Pi_{+}\left(\sigma_{L} \bar{A}\right)\right)-$ $2^{l-1} \cdot 1$, then this exact sequence together with the computation $\operatorname{ind}\left(\rho_{0}-\right.$ $\left.\rho_{1}, \bar{A}\right)=2^{-l}$ shows $\widetilde{K}\left(S\left(T^{*} M\right)\right)$ is generated by $x$ and $y$. We know $2^{l} x=$ 0 and that $\left|\widetilde{K}\left(S\left(T^{*} M\right)\right)\right|=4^{l}$; to determine the additive structure of the group, we must find the order of $y$. Consider the de Rham complex $(d+\delta): C^{\infty}\left(\Lambda^{\text {even }}(D)\right) \rightarrow C^{\infty}\left(\Lambda^{\text {odd }}(D)\right)$ over the disk. The antipodal map acts by +1 on $\Lambda^{\mathrm{e}}$ and by -1 on $\Lambda^{\circ}$. Let $\bar{A}_{1}$ be the tangential operator of this complex. We decompose $(d+\delta)$ into $2^{l}$ operators each of which is isomorphic to $\pm Q$. This indeterminacy does not affect the tangential operator and thus $\bar{A}_{1}=2^{l} \bar{A}$.

The symbol of $\bar{A}_{1}$ on $\Lambda^{\text {even }}(D)$ is $-i c(d n) c(\zeta)$ for $\zeta \in T^{*}(M)$. Let $\tau\left\{\theta_{\text {even }}+\theta_{\text {odd }}\right\}=\theta_{\text {even }}+c(d n) \theta_{\text {odd }}$ provide an isomorphism between $\Lambda(M)$ and $\Lambda^{\text {even }}\left(\left.D\right|_{M}\right)$. We may regard $\bar{A}_{1}$ as an operator on $C^{\infty}(\Lambda(M))$ with symbol $\bar{a}_{1}$ given by:

$$
\begin{aligned}
\bar{a}_{1}(x, \zeta)\left(\theta_{\text {even }}+\theta_{\text {odd }}\right) & =\left\{\tau^{-1} \cdot-i c(d n) c(\zeta) \cdot \tau\right\}\left\{\theta_{\text {even }}+\theta_{\text {odd }}\right\} \\
& =\left\{\tau^{-1} \cdot-i c(d n) c(\zeta)\right\}\left\{\theta_{\text {even }}+c(d n) \theta_{\text {odd }}\right\} \\
& =\tau^{-1}\left\{-i c(d n) c(\zeta) \theta_{\text {even }}-i c(\zeta) \theta_{\text {odd }}\right\} \\
& =-i c(\zeta)\left\{\theta_{\text {even }}+\theta_{\text {odd }}\right\}
\end{aligned}
$$

so that $\bar{A}_{1}=-(d+\delta)$ on $C^{\infty}(\Lambda(M))$. Let $\varepsilon\left(\theta_{\text {odd }}\right)=\theta_{\text {odd }}-i c(\zeta) \theta_{\text {odd }}$ provide an isomorphism between $\Lambda^{\text {odd }}(M)$ and $\Pi_{+}\left(\bar{a}_{1}\right)$. This shows:

$$
2^{l} \cdot y=\Lambda^{\text {odd }}(M)-2^{2 l-1} \cdot 1
$$

Let $\gamma(W)=W-\operatorname{dim}(W) \cdot 1$ be the natural projection of $K(M)$ on $\widetilde{K}(M)$. We wish to compute $\gamma\left(\Lambda^{\text {odd }}(M)\right)$. As complex vector bundles we have $T^{*}(M) \oplus 1=(2 l+1) \cdot L$ so that we have the relation $\Lambda^{j}\left(T^{*} M\right) \oplus$ $\Lambda^{j-1}\left(T^{*} M\right)=\binom{2 l+1}{j} \cdot L^{j}$. This yields the identities:

$$
\begin{array}{ll}
\gamma\left(\Lambda^{j}\left(T^{*} M\right)\right)+\gamma\left(\Lambda^{j-1}\left(T^{*} M\right)\right)=0 & \text { if } j \text { is even } \\
\gamma\left(\Lambda^{j}\left(T^{*} M\right)\right)+\gamma\left(\Lambda^{j-1}\left(T^{*} M\right)\right)=\binom{2 l+1}{j} \cdot \gamma(L) & \text { if } j \text { is odd. }
\end{array}
$$

Thus $\gamma\left(\Lambda^{j}\left(T^{*} M\right)\right)=\gamma\left(\Lambda^{j-2}\left(T^{*} M\right)\right)+\binom{2 l+1}{j} \gamma(L)$ if $j$ is odd. This leads to the identity:

$$
\begin{aligned}
\gamma\left(\Lambda^{2 j+1}\left(T^{*} M\right)\right) & =\left\{\binom{2 l+1}{2 j+1}+\binom{2 l+1}{2 j-1}+\cdots+\binom{2 l+1}{1}\right\} \cdot \gamma(L) \\
\gamma\left(\Lambda^{\text {odd }}\left(T^{*} M\right)\right) & =\left\{l\binom{2 l+1}{1}+(l-1)\binom{2 l+1}{3}+\cdots+\binom{2 l+1}{2 l-1}\right\} \cdot \gamma(L) .
\end{aligned}
$$

We complete this calculation by evaluating this coefficient. Let

$$
\begin{aligned}
f(t) & =\frac{1}{2}\left((t+1)^{2 l+1}-(t-1)^{2 l+1}\right) \\
& =t^{2 l}\binom{2 l+1}{1}+t^{2 l-2}\binom{2 l+1}{3}+\cdots+t^{2}\binom{2 l+1}{2 l+1}+1 \\
f^{\prime}(t) & =\frac{1}{2}(2 l+1)\left((t+1)^{2 l}-(t-1)^{2 l}\right) \\
& =2 \cdot\left\{t^{2 l-1} \cdot l \cdot\binom{2 l+1}{1}+t^{2 l-3} \cdot(l-1) \cdot\binom{2 l+1}{3}+\cdots+t \cdot 1 \cdot\binom{2 l+1}{2 l-1}\right\} .
\end{aligned}
$$

We evaluate at $t=1$ to conclude:

$$
\begin{aligned}
l\binom{2 l+1}{1}+(l-1)\binom{2 l+1}{3}+\cdots+ & \binom{2 l+1}{2 l-1} \\
& =\frac{1}{2} f^{\prime}(1)=\frac{1}{4}(2 l+1) 2^{2 l}=(2 l+1) 4^{l-1}
\end{aligned}
$$

Therefore:

$$
2^{l} \cdot y=(2 l+1) \cdot 4^{l-1}(L-1)
$$

If $l=1$, this gives the relation $2 y=3 x=x$ so $\widetilde{K}\left(S\left(T^{*} M\right)\right)=\mathbf{Z}_{4}$. In fact, $S\left(T^{*} M\right)=S^{3} / \mathbf{Z}_{4}$ is a lens space so this calculation agrees with Corollary 4.6.10. If $l>1$, then $2^{l} \mid 4^{l-1}$ so $2^{l} \cdot y=0$. From this it follows $\widetilde{K}\left(S\left(T^{*} M\right)\right)=\mathbf{Z} / 2^{l} \mathbf{Z} \oplus \mathbf{Z} / 2^{l} \mathbf{Z}$ and the short exact sequence actually splits in this case. This proves:

Theorem 4.6.12 (V. Snaith). Let $X=S\left(T^{*}\left(\mathbf{R} P^{2 l}\right)\right.$ ) be the unit tangent bundle over even dimenional real projective space. If $l=1$ then $K(X)=\mathbf{Z} \oplus \mathbf{Z}_{4}$. Otherwise $K(X)=\mathbf{Z} \oplus \mathbf{Z} / 2^{l} \mathbf{Z} \oplus \mathbf{Z} / 2^{l} \mathbf{Z}$. The map
$\operatorname{ind}(*, *)$ gives a perfect pairing $K\left(\mathbf{R} P^{2 l}\right) \otimes K(X) / K\left(\mathbf{R} P^{2 l}\right) \rightarrow \mathbf{Q} \bmod \mathbf{Z}$. The generators of $K(X)$ are $\{1, x, y\}$ for $x=L-1$ and $y=\gamma \Pi_{+}(\bar{a})$.

Remark: This gives the additive structure. We have $x^{2}=-2 x$. We can calculate $x \cdot y$ geometrically. Let $p(u)=i \sum v_{j} \cdot e_{j}$. We regard $p: 1^{v} \rightarrow L^{v}$ for $v=2^{l}$ over $\mathbf{R} P_{2} l$. Let $\bar{a}_{L}=a \otimes I$ on $L^{v}$. Then $p a_{L} p=\bar{a}_{L}$ so $\Pi_{ \pm}(\bar{a}) \otimes L=\Pi_{\mp}(\bar{a})$. Therefore:

$$
\begin{aligned}
(L-1) \otimes\left(\Pi_{+}(\bar{a})-1^{2^{l-1}}\right) & =-2^{l-1}(L-1)+\left(L \otimes \Pi_{+}(\bar{a})-\Pi_{+}(\bar{a})\right) \\
& =2^{l-1}(L-1)+\Pi_{-}(\bar{a})-\Pi_{+}(\bar{a}) \\
& =2^{l-1}(L-1)+\Pi_{+}(\bar{a})+\Pi_{-}(\bar{a})-2 \Pi_{+}(\bar{a}) \\
& =2^{l-1} \cdot x-2 y
\end{aligned}
$$

so that $x \cdot y=2^{l-1} \cdot x-2 y$. This gives at least part of the ring structure; we do not know a similar simple geometric argument to compute $y \cdot y$.

### 4.7. Singer's Conjecture for the Euler Form.

In this section, we will study a partial converse to the Gauss-Bonnet theorem as well as other index theorems. This will lead to information regarding the higher order terms which appear in the heat equation. In a lecture at M.I.T., I. M. Singer proposed the following question:

Suppose that $P(G)$ is a scalar valued invariant of the metric such that $P(M)=\int_{M} P(G)$ dvol is independent of the metric. Then is there some universal constant $c$ so that $P(M)=c \chi(M)$ ?

Put another way, the Gauss-Bonnet theorem gives a local formula for a topological invariant (the Euler characteristic). Is this the only theorem of its kind? The answer to this question is yes and the result is due to E. Miller who settled the conjecture using topological means. We also settled the question in the affirmative using local geometry independently and we would like to present at least some the the ideas involved. If the invariant is allowed to depend upon the orientation of the manifold, then the characteristic numbers also enter as we shall see later.

We let $P\left(g_{i j / \alpha}\right)$ be a polynomial invariant of the metric; real analytic or smooth invariants can be handled similarly. We suppose $P(M)=$ $\int_{M} P(G)$ dvol is independent of the particular metric chosen on $M$.

Lemma 4.7.1. Let $P$ be a polynomial invariant of the metric tensor with coefficients which depend smoothly on the $g_{i j}$ variables. Suppose $P(M)$ is independent of the metric chosen. We decompose $P=\sum P_{n}$ for $P_{n} \in \mathcal{P}_{m, n}$ homogeneous of order $n$ in the metric. Then $P_{n}(M)$ is independent of the metric chosen separately for each $n ; P_{n}(M)=0$ for $n \neq m$.

Proof: This lets us reduce the questions involved to the homogeneous case. If we replace the metric $G$ by $c^{2} G$ then $P_{n}\left(c^{2} G\right)=c^{-n} P_{n}(G)$ by Lemma 2.4.4. Therefore $\int_{M} P\left(c^{2} G\right) \operatorname{dvol}\left(c^{2} G\right)=\sum_{n} c^{m-n} \int_{M} P_{n}(G)$. Since this is independent of the constant $c, P_{n}(M)=0$ for $n \leq m$ and $P_{m}(M)=P(M)$ which completes the proof.

If $Q$ is 1 -form valued, we let $P=\operatorname{div} Q$ be scalar valued. It is clear that $\int_{M} \operatorname{div} Q(G) \operatorname{dvol}(G)=0$ so $P(M)=0$ in this case. The following gives a partial converse:

Lemma 4.7.2. Let $P \in \mathcal{P}_{m, n}$ for $n \neq m$ satisfy $P(M)=\int_{M} P(G)$ dvol is independent of the metric $G$. Then there exists $Q \in \mathcal{P}_{m, n-1,1}$ so that $P=\operatorname{div} Q$.

Proof: Since $n \neq m, \quad P(M)=0$. Let $f(x)$ be a real valued function on $M$ and let $G_{t}$ be the metric $e^{t f(x)} G$. If $n=0$, then $P$ is constant so $P=0$ and the lemma is immediate. We assume $n>0$ henceforth. We let
$P(t)=P\left(e^{t f(x)} G\right)$ and compute:

$$
\begin{aligned}
\frac{d}{d t}\left(P\left(e^{t f(x)} G\right) \mathrm{dvol}\left(e^{t f(x)} G\right)\right) & =\frac{d}{d t}\left(P\left(e^{t f(x)} G\right) e^{t m f(x) / 2}\right) \mathrm{d} \operatorname{vol}(G) \\
& =Q(G, f) \mathrm{dvol}(G)
\end{aligned}
$$

$Q(f, G)$ is a certain expression which is linear in the derivatives of the scaling function $f$. We let $f_{; i_{1} \ldots i_{p}}$ denote the multiple covariant derivatives of $f$ and let $f_{: i_{1} \ldots i_{p}}$ denote the symmetrized covariant derivatives of $f$. We can express

$$
Q(f, G)=\sum Q_{i_{1} \ldots i_{p}} f_{: i_{1} \ldots i_{p}}
$$

where the sum ranges over symmetric tensors of length less than $n$. We formally integrate by parts to express:

$$
Q(f, G)=\operatorname{div} R(f, G)+\sum(-1)^{p} Q_{i_{1} \ldots i_{p}: i_{1} \ldots i_{p}} f .
$$

By integrating over the structure group $\mathrm{O}(m)$ we can ensure that this process is invariantly defined. If $S(G)=\sum(-1)^{p} Q_{i_{1} \ldots i_{p}: i_{1} \ldots i_{p}}$ then the identity:

$$
0=\int_{M} Q(f, G) \operatorname{dvol}(G)=\int_{M} S(G) f \mathrm{dvol}(G)
$$

is true for every real valued function $f$. This implies $S(G)=0$ so $Q(f, G)=$ $\operatorname{div} R(f, G)$. We set $f=1$. Since

$$
e^{t m / 2} P\left(e^{t} G\right)=e^{(m-n) t / 2} P(G),
$$

we conclude $Q(1, G)=\frac{m-n}{2} P(G)$ so $P(G)=\frac{2}{m-n} \operatorname{div} R(1, G)$ which completes the proof.

There is a corresponding lemma for form valued invariants. The proof is somewhat more complicated and we refer to (Gilkey, Smooth invariants of a Riemannian manifold) for further details:

Lemma 4.7.3.
(a) Let $P \in \mathcal{P}_{m, n, p}$. We assume $n \neq p$ and $d P=0$. If $p=m$, we assume $\int_{M} P(G)$ is independent of $G$ for every $G$ on $M$. Then there exists $Q \in \mathcal{P}_{m, n-1, p-1}$ so that $d Q=P$.
(b) Let $P \in \mathcal{P}_{m, n, p}$. We assume $n \neq m-p$ and $\delta P=0$. If $p=0$ we assume $\int_{M} P(G)$ dvol is independent of $G$ for every $G$ in $M$. Then there exists $Q \in \mathcal{P}_{m, n-1, p+1}$ so that $\delta Q=P$.

Remark: (a) and (b) are in a sense dual if one works with $\mathrm{SO}(m)$ invariance and not just $O(m)$ invariance. We can use this Lemma together with the results of Chapter 2 to answer the generalized Singer's conjecture:

Theorem 4.7.4.
(a) Let $P$ be a scalar valued invariant so that $P(M)=\int_{M} P(G) \mathrm{dvol}$ is independent of the particular metric chosen. Then we can decompose $P=$ $c \cdot E_{m}+\operatorname{div} Q$ where $E_{m}$ is the Euler form and where $Q$ is a 1-form valued invariant. This implies $P(M)=c \chi(M)$.
(b) Let $P$ be a $p$-form valued invariant so that $d P=0$. If $p=m$, we assume $P(M)=\int_{M} P(G)$ is independent of the particular metric chosen. Then we can decompose $P=R+d Q . \quad Q$ is $p-1$ form valued and $R$ is a Pontrjagin form.

Proof: We decompose $P=\sum P_{j}$ into terms which are homogeneous of order $j$. Then each $P_{j}$ satisfies the hypothesis of Theorem 4.7.4 separately so we may assume without loss of generality that $P$ is homogeneous of order $n$. Let $P$ be as in (a). If $n \neq m$, then $P=\operatorname{div} Q$ be Lemma 4.7.2. If $n=m$, we let $P_{1}=r(P) \in \mathcal{P}_{m-1, m}$. It is immediate $\int_{M_{1}} P_{1}\left(G_{1}\right) \operatorname{dvol}\left(G_{1}\right)=$ $\frac{1}{2 \pi} \int_{S^{1} \times M_{1}} P\left(1 \times G_{1}\right)$ is independent of the metric $G_{1}$ so $P_{1}$ satisfies the hypothesis of (a) as well. Since $m-1 \neq m$ we conclude $P_{1}=\operatorname{div} Q_{1}$ for $Q_{1} \in \mathcal{P}_{m-1, m-1,1}$. Since $r$ is surjective, we can choose $Q \in \mathcal{P}_{m, m-1,1}$ so $r(Q)=Q_{1}$. Therefore $r(P-\operatorname{div} Q)=P_{1}-\operatorname{div} Q_{1}=0$ so by Theorem 2.4.7, $P-\operatorname{div} Q=c E_{m}$ for some constant $c$ which completes the proof. (The fact that $r: \mathcal{P}_{m, n, p} \rightarrow \mathcal{P}_{m-1, n, p} \rightarrow 0$ is, of course, a consequence of H. Weyl's theorem so we are using the full force of this theorem at this point for the first time).

If $P$ is $p$-form valued, the proof is even easier. We decompose $P$ into homogeneous terms and observe each term satisfies the hypothesis separately. If $P$ is homogeneous of degree $n \neq p$ then $P=d Q$ be Lemma 4.7.3. If $P$ is homogeneous of degree $n=p$, then $P$ is a Pontrjagin form by Lemma 2.5.6 which completes the proof in this case.

The situation in the complex catagory is not as satisfactory.
Theorem 4.7.5. Let $M$ be a holomorphic manifold and let $P$ be a scalar valued invariant of the metric. Assume $P(M)=\int_{M} P(G)$ dvol is independent of the metric $G$. Then we can express $P=R+\operatorname{div} Q+\mathcal{E} . Q$ is a 1 -form valued invariant and $P=* R^{\prime}$ where $R^{\prime}$ is a Chern form. The additional error term $\mathcal{E}$ satisfies the conditions: $r(\mathcal{E})=0$ and $\mathcal{E}$ vanishes for Kaehler metric. Therefore $P(M)$ is a characteristic number if $M$ admits a Kaehler metric.

The additional error term arises because the axiomatic characterization of the Chern forms given in Chapter 3 were only valid for Kaehler metrics. $\mathcal{E}$ in general involves the torsion of the holomorphic connection and to show $\operatorname{div} Q^{\prime}=\mathcal{E}$ for some $Q^{\prime}$ is an open problem. Using the work of E. Miller, it does follow that $\int \mathcal{E}$ dvol $=0$ but the situation is not yet completely resolved.

We can use these results to obtain further information regarding the higher terms in the heat expansion:

Theorem 4.7.6.
(a) Let $a_{n}(x, d+\delta)$ denote the invariants of the heat equation for the de Rham complex. Then
(i) $a_{n}(x, d+\delta)=0$ if $m$ or $n$ is odd or if $n<m$,
(ii) $a_{m}(x, d+\delta)=E_{m}$ is the Euler form,
(iii) If $m$ is even and if $n>m$, then $a_{n}(x, d+\delta) \not \equiv 0$ in general. However, there does exist a 1 -form valued invariant $q_{m, n}$ so $a_{n}=\operatorname{div} q_{m, n}$ and $r\left(q_{m, n}\right)=0$.
(b) Let $a_{n}^{\text {sign }}(x, V)$ be the invariants of the heat equation for the signature complex with coefficients in a bundle $V$. Then
(i) $a_{n}^{\text {sign }}(x, V)=0$ for $n<m$ or $n$ odd,
(ii) $a_{m}^{\text {sign }}(x, V)$ is the integrand of the Atiyah-Singer index theorem,
(iii) $a_{n}^{\text {sign }}(x, V)=0$ for $n$ even and $n<m$. However, there exists an $m-1$ form valued invariant $q_{m, n}^{\text {sign }}(x, V)$ so that $a_{n}^{\text {sign }}=d\left(q_{m, n}^{\text {sign }}\right)$.
A similar result holds for the invariants of the spin complex.
(c) Let $a_{n}^{\text {Dolbeault }}(x, V)$ denote the invariants of the heat equation for the Dolbeault complex with coefficients in $V$. We do not assume that the metric in question is Kaehler. Then:
(i) $a_{n}^{\text {Dolbeault }}(x, V)=\operatorname{div} Q_{m, n}$ where $Q_{m, n}$ is a 1-form valued invariant for $n \neq m$,
(ii) $a_{m}^{\text {Dolbeault }}(x, V)=\operatorname{div} Q_{m, n}+$ the integrand of the Riemann-Roch theorem.

The proof of all these results relies heavily on Lemma 4.7.3 and 4.7.2 and will be omitted. The results for the Dolbeault complex are somewhat more difficult to obtain and involve a use of the $\mathrm{SPIN}_{c}$ complex.

We return to the study of the de Rham complex. These arguments are due to L. Willis. Let $m=2 n$ and let $a_{m+2}(x, d+\delta)$ be the next term above the Euler form. Then $a_{m+2}=\operatorname{div} Q_{m+1}$ where $Q_{m+1} \in \mathcal{P}_{m, m+1,1}$ satisfies $r\left(Q_{m+1}\right)=0$. We wish to compute $a_{m+2}$. The first step is:

Lemma 4.7.7. Let $m$ be even and let $Q \in \mathcal{P}_{m, m+1,1}$ be 1 -form valued. Suppose $r(Q)=0$. Then $Q$ is a linear combination of $d E_{m}$ and $\Phi_{m}$ defined by:

$$
\begin{aligned}
& \Phi_{m}=\sum_{k, \rho, \tau} \operatorname{sign}(\rho) \operatorname{sign}(\tau) \cdot\left\{(-8 \pi)^{m / 2}\left(\frac{m}{2}-1\right)!\right\}^{-1} \\
& \times R_{\rho(1) \rho(2) \tau(1) k ; k} R_{\rho(3) \rho(4) \tau(3) \tau(4)} \cdots R_{\rho(m-1) \rho(m) \tau(m-1) \tau(m)} e^{\tau(2)} \in \Lambda^{1}
\end{aligned}
$$

where $\left\{e^{1}, \ldots, e^{m}\right\}$ are a local orthonormal frame for $T^{*} M$.

Proof: We let $A e^{k}$ be a monomial of $P$. We express $A$ in the form:

$$
A=g_{i_{1} j_{1} / \alpha_{1}} \ldots g_{i_{r} j_{r} / \alpha_{r}}
$$

By Lemma 2.5.5, we could choose a monomial $B$ of $P$ with $\operatorname{deg}_{k}(B)=0$ for $k>2 \ell(A)=2 r$. Since $r(Q)=0$, we conclude $2 r \geq m$. On the other hand, $2 r \leq \sum\left|\alpha_{\nu}\right|=m+1$ so we see

$$
m \leq 2 r \leq m+1
$$

Since $m$ is even, we conclude $2 r=m$. This implies one of the $\left|\alpha_{\nu}\right|=3$ while all the other $\left|\alpha_{\nu}\right|=2$. We choose the notation so $\left|\alpha_{1}\right|=3$ and $\left|\alpha_{\nu}\right|=2$ for $\nu>1$. By Lemma 2.5.1, we may choose $A$ in the form:

$$
A=g_{i j / 111} g_{i_{2} j_{2} / k_{2} l_{2}} \ldots g_{i_{r} j_{r} / k_{r} l_{r}}
$$

We first suppose $\operatorname{deg}_{1}(A)=3$. This implies $A$ appears in the expression $A e^{1}$ in $P$ so $\operatorname{deg}_{j}(A) \geq 2$ is even for $j>1$. We estimate $2 m+1=$ $4 r+1=\sum_{j} \operatorname{deg}_{j}(A)=3+\sum_{j>1} \operatorname{deg}_{j}(A) \geq 3+2(m-1)=2 m+1$ to show $\operatorname{deg}_{j}(A)=2$ for $j>1$. If $m=2$, then $A=g_{22 / 111}$ which shows $\operatorname{dim} \mathrm{N}(r)=1$ and $Q$ is a multiple of $d E_{2}$. We therefore assume $m>2$. Since $\operatorname{deg}_{j}(A)=2$ for $j>1$, we can apply the arguments used to prove Theorem 2.4.7 to the indices $j>1$ to show that

$$
g_{22 / 111} g_{33 / 44} \ldots g_{m-1 / m m} e^{1}
$$

is a monomial of $P$.
Next we suppose $\operatorname{deg}_{1}(A)>3$. If $\operatorname{deg}_{1}(A)$ is odd, then $\operatorname{deg}_{j}(A) \geq 2$ is even for $j \geq 2$ implies $2 m+1=4 r+1=\sum_{j} \operatorname{deg}_{j}(A) \geq 5+2(m-1)=2 m+3$ which is false. Therefore $\operatorname{deg}_{1}(A)$ is even. We choose the notation in this case so $A e^{2}$ appears in $P$ and therefore $\operatorname{deg}_{2}(A)$ is odd and $\operatorname{deg}_{j}(A) \geq 2$ for $j>2$. This implies $2 m+1=4 r+1=\operatorname{deg}_{j}(A) \geq 4+1+\sum_{j>2} \operatorname{deg}_{j}(A) \geq$ $5+2(m-2)=2 m+1$. Since all the inequalities must be equalities we conclude

$$
\operatorname{deg}_{1}(A)=4, \quad \operatorname{deg}_{2}(A)=1, \quad \operatorname{deg}_{3}(A)=2 \text { for } j>2
$$

We apply the arguments of the second chapter to choose $A$ of this form so that every index $j>2$ which does not touch either the index 1 or the index 2 touches itself. We choose $A$ so the number of indices which touch themselves in $A$ is maximal. Suppose the index 2 touches some other index than the index 1. If the index 2 touches the index 3 , then the index 3 cannot touch itself in $A$. An argument using the fact $\operatorname{deg}_{2}(A)=1$ and using the arguments of the second chapter shows this would contradict the maximality of $A$ and thus the index 1 must touch the index 2 in $A$. We use
non-linear changes of coordinates to show $g_{12 / 111}$ cannot divide $A$. Using non-linear changes of coordinates to raplace $g_{44 / 12}$ by $g_{12 / 44}$ if necessary, we conclude $A$ has the form

$$
A=g_{33 / 111} g_{12 / 44} g_{55 / 66} \ldots g_{m-1, m-1 / m m}
$$

This shows that if $m \geq 4$, then $\operatorname{dim} \mathrm{N}\left(r: \mathcal{P}_{m, m+1,1} \rightarrow \mathcal{P}_{m-1, m+1,1}\right) \leq 2$. Since these two invariants are linearly independent for $m \geq 4$ this completes the proof of the lemma.

We apply the lemma and decompose

$$
a_{m+2}(x, d+\delta)=c_{1}(m)\left(E_{m}\right)_{; k k}+c_{2}(m) \operatorname{div} \Phi_{m}
$$

The $c_{1}(m)$ and $c_{2}(m)$ are universal constants defending only on the dimension $m$. We consider a product manifold of the form $M=S^{2} \times M^{m-2}$ to define a map:

$$
r_{(2)}: \mathcal{P}_{m, n, 0} \rightarrow \bigoplus_{q \leq n} \mathcal{P}_{m-2, q, 0}
$$

This restriction map does not preserve the grading since by throwing derivatives to the metric over $S^{2}$, we can lower the order of the invariant involved.

Let $x_{1} \in S^{2}$ and $x_{2} \in M^{m-2}$. The multiplicative nature of the de Rham complex implies:

$$
a_{n}(x, d+\delta)=\sum_{j+k=m} a_{j}\left(x_{1}, d+\delta\right)_{S^{2}} a_{k}\left(x_{2}, d+\delta\right)_{M^{m-2}}
$$

However, $a_{j}\left(x_{1}, d+\delta\right)$ is a constant since $S^{2}$ is a homogeneous space. The relations $\int_{S^{2}} a_{j}\left(x_{1}, d+\delta\right)=2 \delta_{j, 2}$ implies therefore

$$
a_{n}(x, d+\delta)=\frac{1}{2 \pi} a_{n-2}\left(x_{2}, d+\delta\right)
$$

so that

$$
r_{(2)} a_{m+2}^{m}=\frac{1}{2 \pi} a_{m}^{m-2}
$$

Since $r_{(2)}\left(E_{m}\right)_{; k k}=\frac{1}{2 \pi}\left(E_{m-2}\right)_{; k k}$ and $r_{(2)} \operatorname{div} \Phi_{m}=\frac{1}{2 \pi} \operatorname{div} \Phi_{m-2}$, we conclude that in fact the universal constants $c_{1}$ and $c_{2}$ do not depend upon $m$. (If $m=2$, these two invariants are not linearly independent so we adjust $c_{1}(2)=c_{1}$ and $\left.c_{2}(2)=c_{2}\right)$. It is not difficult to use Theorem 4.8.16 which will be discussed in the next section to compute that if $m=4$, then:

$$
a_{6}(x, d+\delta)=\frac{1}{12}\left(E_{m}\right)_{; k k}-\frac{1}{6} \operatorname{div} \Phi_{m} .
$$

We omit the details of the verification. This proves:

Theorem 4.7.8. Let $m$ be even and let $a_{m+2}(x, d+\delta)$ denote the invariant of the Heat equation for the de Rham complex. Then:

$$
a_{m+2}=\frac{1}{12}\left(E_{m}\right)_{; k k}-\frac{1}{6} \operatorname{div} \Phi_{m}
$$

We gave a different proof of this result in (Gilkey, Curvature and the heat equation for the de Rham complex). This proof, due to Willis, is somewhat simpler in that it uses Singer's conjecture to simplify the invariance theory involved.

There are many other results concerning the invariants which appear in the heat equation for the de Rham complex. Gunther and Schimming have given various shuffle formulas which generalize the alternating sum defined previously. The combinatorial complexities are somewhat involved so we shall simply give an example of the formulas which can be derived.

Theorem 4.7.9. Let $m$ be odd. Then

$$
\sum(-1)^{p}(m-p) a_{n}\left(x, \Delta_{p}^{m}\right)= \begin{cases}0 & \text { if } n<m-1 \\ E_{m-1} & \text { if } n=m-1\end{cases}
$$

Proof: We remark that there are similar formulas giving the various Killing curvatures $E_{k}$ for $k<m$ in all dimensions. Since $r: \mathcal{P}_{m, n} \rightarrow \mathcal{P}_{m-1, n}$ for $n<m$ is an isomorphism, it suffices to prove this formula under restriction. Since $r\left(a_{n}\left(x, \Delta_{p}^{m}\right)\right)=a_{n}\left(x, \Delta_{p}^{m-1}\right)+a_{n}\left(x, \Delta_{p-1}^{m-1}\right)$ we must study

$$
\sum(-1)^{p}(m-p)\left(a_{n}\left(x, \Delta_{p}^{m-1}\right)+a_{n}\left(x, \Delta_{p-1}^{m-1}\right)\right)=(-1)^{p} a_{n}\left(x, \Delta_{p}^{m-1}\right)
$$

and apply Theorem 2.4.8.
We remark that all the shuffle formulas of Gunther and Schimming (including Theorem 4.7.9) have natural analogues for the Dolbeault complex for a Kaehler metric and the proofs are essentially the same and rely on Theorem 3.6.9 and 3.4.10.

### 4.8. Local Formulas for the Invariants <br> Of the Heat Equation.

In this subsection, we will compute $a_{n}\left(x, \Delta_{p}^{m}\right)$ for $n=0,1,2$. In principle, the combinatorial formulas from the first chapter could be used in this calculation. In practice, however, these formulas rapidly become much too complicated for practical use so we shall use instead some of the functorial properties of the invariants involved.

Let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be a second order elliptic operator with leading symbol given by the metric tensor. This means we can express:

$$
P=-\left(\sum_{i, j} g^{i j} \partial^{2} / \partial x_{i} \partial x_{j}+\sum_{k} A_{k} \partial / \partial x_{k}+B\right)
$$

where the $A_{k}, B$ are endomorphisms of the bundle $V$ and where the leading term is scalar. This category of differential operators includes all the Laplacians we have been considering previously. Our first task is to get a more invariant formulation for such an operator.

Let $\nabla$ be a connection on $V$ and let $E \in C^{\infty}(\operatorname{END}(V))$ be an endomorphism of $V$. We use the Levi-Civita connection on $M$ and the connection $\nabla$ on $V$ to extend $\nabla$ to tensors of all orders. We define the differential operator $P_{\nabla}$ by the diagram:

$$
P_{\nabla}: C^{\infty}(V) \xrightarrow{\nabla} C^{\infty}\left(T^{*} M \otimes V\right) \xrightarrow{\nabla} C^{\infty}\left(T^{*} M \otimes T^{*} M \otimes V\right) \xrightarrow{-g \otimes 1} C^{\infty}(V) .
$$

Relative to a local orthonormal frame for $T^{*} M$, we can express

$$
P_{\nabla}(f)=-f_{; i i}
$$

so this is the trace of second covariant differentiation. We define

$$
P(\nabla, E)=P_{\nabla}-E
$$

and our first result is:
Lemma 4.8.1. Let $P: C^{\infty}(V) \rightarrow C^{\infty}(V)$ be a second order operator with leading symbol given by the metric tensor. There exists a unique connection on $V$ and a unique endomorphism so $P=P(\nabla, E)$.

Proof: We let indices $i, j, k$ index a coordinate frame for $T(M)$. We shall not introduce indices to index a frame for $V$ but shall always use matrix notation. The Christoffel symbols $\Gamma_{i j}{ }^{k}=-\Gamma_{i}{ }^{k}{ }_{j}$ of the Levi-Civita connection are given by:

$$
\begin{aligned}
\Gamma_{i j}^{k} & =\frac{1}{2} g^{k l}\left(g_{i l / j}+g_{j l / i}-g_{i j / l}\right) \\
\nabla_{\partial / \partial x_{i}}\left(\partial / \partial x_{j}\right) & =\Gamma_{i j}^{k}\left(\partial / \partial x_{k}\right) \\
\nabla_{\partial / \partial x_{i}}\left(d x_{j}\right) & =\Gamma_{i}^{j}{ }_{k}\left(d x_{k}\right)
\end{aligned}
$$

where we sum over repeated indices. We let $\vec{s}$ be a local frame for $V$ and let $\omega_{i}$ be the components of the connection 1-form $\omega$ so that:

$$
\nabla(f \cdot \vec{s})=d x^{i} \otimes\left(\partial f / \partial x_{i}+\omega_{i}(f)\right) \cdot \vec{s}
$$

With this notational convention, the curvature is given by:

$$
\Omega_{i j}=\omega_{j / i}-\omega_{i / j}+\omega_{i} \omega_{j}-\omega_{j} \omega_{i} \quad(\Omega=d \omega+\omega \wedge \omega)
$$

We now compute:

$$
\begin{aligned}
\nabla^{2}(f \cdot \vec{s})=d x^{j} \otimes d x^{i} \otimes\left(\partial^{2} f\right. & \partial x_{i} \partial x_{j}+\omega_{j} \partial f / \partial x_{i}+\Gamma_{j}{ }^{k} \partial f / \partial x_{k} \\
& \left.+\omega_{i} \partial f / \partial x_{j}+\omega_{i / j} f+\omega_{j} \omega_{i} f+\Gamma_{j}{ }^{k}{ }_{i} \omega_{k} f\right) \cdot \vec{s}
\end{aligned}
$$

from which it follows that:

$$
\begin{array}{r}
P_{\nabla}(f \cdot \vec{s})=-\left\{g^{i j} \partial f / \partial x_{i} \partial x_{j}+\left(2 g^{i j} \omega_{i}-g^{j k} \Gamma_{j k}{ }^{i}\right) \partial f / \partial x_{i}\right. \\
\left.+\left(g^{i j} \omega_{i / j}+g^{i j} \omega_{i} \omega_{j}-g^{i j} \Gamma_{i j}{ }^{k} \omega_{k}\right) f\right\} \cdot \vec{s}
\end{array}
$$

We use this identity to compute:

$$
\left(P-P_{\nabla}\right)(f \cdot \vec{s})=-\left\{\left(A_{i}-2 g^{i j} \omega_{j}+g^{j k} \Gamma_{j k}^{i}\right) \partial f / \partial x_{i}+(*) f\right\} \cdot \vec{s}
$$

where we have omitted the $0^{\text {th }}$ order terms. Therefore $\left(P-P_{\nabla}\right)$ is a $0^{\text {th }}$ order operator if and only if

$$
A_{i}-2 g^{i j} \omega_{j}+g^{j k} \Gamma_{j k}^{i}=0 \quad \text { or equivalently } \omega_{i}=\frac{1}{2}\left(g_{i j} A_{j}+g_{i j} g^{k l} \Gamma_{k l}^{j}\right)
$$

This shows that the $\left\{\omega_{i}\right\}$ are uniquely determined by the condition that $\left(P-P_{\nabla}\right)$ is a $0^{\text {th }}$ order operator and specificies the connection $\nabla$ uniquely.

We define:

$$
E=P_{\nabla}-P \text { so } E=B-g^{i j} \omega_{i / j}-g^{i j} \omega_{i} \omega_{j}+g^{i j} \omega_{k} \Gamma_{i j}^{k}
$$

We fix this connection $\nabla$ and endomorphism $E$ determined by $P$. We summarize these formulas as follows:

Corollary 4.8.2. If $(\nabla, E)$ are determined by the second order operator $P=-\left(g^{i j} \partial^{2} f / \partial x_{i} \partial x_{j}+A_{i} \partial j / \partial x_{i}+B f\right) \cdot \vec{s}$, then

$$
\begin{aligned}
\omega_{i} & =\frac{1}{2}\left(g_{i j} A_{j}+g_{i j} g^{k l} \Gamma_{k l}^{j}\right) \\
E & =B-g^{i j} \omega_{i / j}-g^{i j} \omega_{i} \omega_{j}+g^{i j} \omega_{k} \Gamma_{i j}^{k}
\end{aligned}
$$

We digress briefly to express the Laplacian $\Delta_{p}$ in this form. If $\nabla$ is the Levi-Civita connection acting on $p$-forms, it is clear that $\Delta_{p}-P_{\nabla}$ is a first
order operator whose leading symbol is linear in the 1-jets of the metric. Since it is invariant, the leading symbol vanishes so $\Delta_{p}-P_{\nabla}$ is a $0^{\text {th }}$ order and the Levi-Civita connection is in fact the connection determined by the operator $\Delta_{p}$. We must now compute the curvature term. The operator $(d+\delta)$ is defined by Clifford multiplication:

$$
\begin{aligned}
& (d+\delta): C^{\infty}\left(\Lambda\left(T^{*} M\right)\right) \rightarrow C^{\infty}\left(T^{*} M \otimes \Lambda\left(T^{*} M\right)\right) \\
& \xrightarrow{\text { Clifford multiplication }} C^{\infty}\left(\Lambda\left(T^{*} M\right)\right) .
\end{aligned}
$$

If we expand $\theta=f_{I} d x_{I}$ and $\nabla \theta=d x_{i} \otimes\left(\partial f_{I} / \partial x_{i}+\Gamma_{i I}{ }^{J} f_{I}\right) d x_{J}$ then it is immediate that if "*" denotes Clifford multiplication,

$$
\begin{aligned}
(d+\delta)\left(f_{I} d x_{I}\right) & =\left(f_{I / i} d x_{i} * d x_{I}\right)+\left(f_{I} \Gamma_{i I J} d x_{i} * d x_{J}\right) \\
\Delta\left(f_{I} d x_{I}\right) & =\left(f_{I / i j} d x_{j} * d x_{i} * d x_{I}\right)+\left(f_{I} \Gamma_{i I J / j} d x_{j} * d x_{i} * d x_{J}\right)+\cdots
\end{aligned}
$$

where we have omitted terms involving the 1-jets of the metric. Similary, we compute:

$$
P_{\nabla}\left(f_{I} d x_{I}\right)=-g^{i j} f_{I / i j} d x_{I}-g^{i j} \Gamma_{i I J / j} f_{I} d x_{J}+\cdots
$$

We now fix a point $x_{0}$ of $M$ and let $X$ be a system of geodesic polar coordinates centered at $x_{0}$. Then $\Gamma_{i I J / i}=0$ and $\Gamma_{i I J / j}=\frac{1}{2} R_{j i I J}$ gives the curvature tensor at $x_{0}$. Using the identities $d x_{i} * d x_{j}+d x_{j} * d x_{i}=-2 \delta_{i j}$ we see $f_{I / i j} d x_{j} * d x_{i}=-f_{I / i i}$ and consequently:
$\left(P_{\nabla}-\Delta\right)\left(f_{I} d x_{I}\right)=-\frac{1}{2} R_{i j I J} f_{I} d x_{i} * d x_{j} * d x_{J}=\sum_{i<j} R_{i j I J} f_{I} d x_{j} * d x_{i} * d x_{J}$.
This identity holds true at $x_{0}$ in geodesic polar coordinates. Since both sides are tensorial, it holds in general which proves:

Lemma 4.8.3. Let $\Delta_{p}$ be the Laplacian acting on $p$-forms and let $R_{i j I J}$ be the curvature tensor of the Levi-Civita connection. Then

$$
\left(P_{\nabla}-\Delta_{p}\right)\left(f_{I} d x_{I}\right)=\sum_{i<j} R_{i j I J} f_{I} d x_{j} * d x_{i} * d x_{J}
$$

We now return to the problem of computing the invariants $a_{n}(x, P)$.
Lemma 4.8.4. Let $P$ be an operator as in Lemma 4.8.1. Then $a_{0}(x, P)=$ $(4 \pi)^{-m / 2} \operatorname{dim} V$.
Proof: We first consider the operator $P=-\partial^{2} / \partial \theta^{2}$ on the unit circle $[0,2 \pi]$. The eigenvalues of $P$ are $\left\{n^{2}\right\}_{n \in \mathbf{Z}}$. Since $a_{0}$ is homogeneous of
order 0 in the jets of the symbol, $a_{0}(\theta, P)=a_{0}(P) / \operatorname{vol}(M)$ is constant. We compute:

$$
\sum_{n} e^{-t n^{2}} \sim t^{-1 / 2} \int_{M} a_{0}(\theta, P) d \theta=t^{-1 / 2}(2 \pi) a_{0}(\theta, P)
$$

However, the Riemann sums approximating the integral show that

$$
\sqrt{\pi}=\int_{-\infty}^{\infty} e^{-t x^{2}} d x=\lim _{t \rightarrow 0} \sum_{n} e^{-(\sqrt{t} n)^{2}} \sqrt{t}
$$

from which it follows immediately that $a_{0}(\theta, P)=\sqrt{\pi} /(2 \pi)=(4 \pi)^{-1 / 2}$. More generally, by taking the direct sum of these operators acting on $C^{\infty}\left(S^{1} \times \mathbf{C}^{k}\right)$ we conclude $a_{0}(\theta, P)=(4 \pi)^{-1 / 2} \operatorname{dim} V$ which completes the proof if $m=1$.

More generally, $a_{0}(x, P)$ is homogeneous of order 0 in the jets of the symbol so $a_{0}(x, P)$ is a constant which only depends on the dimension of the manifold and the dimension of the vector bundle. Using the additivity of Lemma 1.7.5, we conclude $a_{0}(x, P)$ must have the form:

$$
a_{0}(x, P)=c(m) \cdot \operatorname{dim} V .
$$

We now let $M=S^{1} \times \cdots \times S^{1}$, the flat $m$-torus, and let $\Delta=-\sum_{\nu} \partial^{2} / \partial \theta_{\nu}^{2}$. The product formula of Lemma 1.7.5(b) implies that:

$$
a_{0}(x, \Delta)=\prod_{\nu} a_{0}\left(\theta_{\nu},-\partial^{2} / \partial \theta_{\nu}^{2}\right)=(4 \pi)^{-m / 2}
$$

which completes the proof of the lemma.
The functorial properties of Lemma 1.7.5 were essential to the proof of Lemma 4.8.4. We will continue to exploit this functoriality in computing $a_{2}$ and $a_{4}$. (In principal one could also compute $a_{6}$ in this way, but the calculations become of formidable difficulty and will be omitted).

It is convenient to work with more tensorial objects than with the jets of the symbol of $P$. We let $\operatorname{dim} V=k$ and $\operatorname{dim} M=m$. We introduce formal variables $\left\{R_{i_{1} i_{2} i_{3} i_{4} ; \ldots}, \Omega_{i_{1} i_{2} ; \ldots}, E_{;} \ldots\right\}$ for the covariant derivatives of the curvature tensor of the Levi-Civita connection, of the curvature tensor of $\nabla$, and of the covariant derivatives of the endomorphism $E$. We let $\mathcal{S}$ be the non-commutative algebra generated by these variables. Since there are relations among these variables, $\mathcal{S}$ isn't free. If $S \in \mathcal{S}$ and if $e$ is a local orthonormal frame for $T^{*}(M)$, we define $S(P)(e)(x)=$ $S(G, \nabla, E)(e)(x)(\operatorname{END}(V))$ by evaluation. We say $S$ is invariant if $S(P)=$ $S(P)(a)$ is independent of the orthonormal frame $e$ chosen for $T(M)$. We define:

$$
\begin{aligned}
\operatorname{ord}\left(R_{i_{1} i_{2} i_{4} i_{4} ; i_{5} \ldots i_{k}}\right) & =k+2 \\
\operatorname{ord}\left(\Omega_{i_{1} i_{2} ; i_{3} \ldots i_{k}}\right) & =k \\
\operatorname{ord}\left(E_{; i_{1} \ldots i_{k}}\right) & =k+2
\end{aligned}
$$

to be the degree of homogeneity in the jets of the symbol of $P$, and let $\mathcal{S}_{m, n, k}$ be the finite dimensional subspace of S consisting of the invariant polynomials which are homogeneous of order $n$.

If we apply H . Weyl's theorem to this situation and apply the symmetries involved, it is not difficult to show:

## Lemma 4.8.5.

(a) $\mathcal{S}_{m, 2, k}$ is spanned by the two polynomials $R_{i j i j} I, E$.
(b) $\mathcal{S}_{m, 4, k}$ is spanned by the eight polynomials

$$
\begin{array}{cccc}
R_{i j i j ; k k} I, & R_{i j i j} R_{k l k l} I, & R_{i j i k} R_{l j l k} I, & R_{i j k l} R_{i j k l} I, \\
E^{2}, & E_{; i i}, & R_{i j i j} E, & \Omega_{i j} \Omega_{i j} .
\end{array}
$$

We omit the proof in the interests of brevity; the corresponding spanning set for $\mathcal{S}_{m, 6, k}$ involves 46 polynomials.

The spaces $\mathcal{S}_{m, n, k}$ are related to the invariants $a_{n}(x, P)$ of the heat equation as follows:

Lemma 4.8.6. Let $(m, n, k)$ be given and let $P$ satisfy the hypothesis of Lemma 4.8.1. Then there exists $S_{m, n, k} \in \mathcal{S}_{m, n, k}$ so that $a_{n}(x, P)=$ $\operatorname{Tr}\left(S_{m, n, k}\right)$.

Proof: We fix a point $x_{0} \in M$. We choose geodesic polar coordinates centred at $x_{0}$. In such a coordinate system, all the jets of the metric at $x_{0}$ can be computed in terms of the $R_{i_{1} i_{2} i_{3} i_{4} ; \ldots}$ variables. We fix a frame $s_{0}$ for the fiber $V_{0}$ over $x_{0}$ and extend $s_{0}$ by parallel translation along all the geodesic rays from $x_{0}$ to get a frame near $x_{0}$. Then all the derivatives of the connection 1-form at $x_{0}$ can be expressed in terms of the $R_{\ldots}$.. and $\Omega_{\ldots}$.. variables at $x_{0}$. We can solve the relations of Corollary 4.8.2 to express the jets of the symbol of $P$ in terms of the jets of the metric, the jets of the connection 1-form, and the jets of the endomorphism $E$. These jets can all be expressed in terms of the variables in $\mathcal{S}$ at $x_{0}$ so any invariant endomorphism which is homogeneous of order $n$ belongs to $\mathcal{S}_{m, n, k}$. In Lemma 1.7.5, we showed that $a_{n}(x, P)=\operatorname{Tr}\left(e_{n}(x, P)\right)$ was the trace of an invariant endomorphism and this completes the proof; we set $S_{m, n, k}=e_{n}(x, P)$.

We use Lemmas 4.8.5 and 4.8.6 to expand $a_{n}(x, P)$. We regard scalar invariants of the metric as acting on $V$ by scalar multiplication; alternatively, such an invariant $R_{i j i j}$ could be replaced by $R_{i j i j} I_{V}$.

Lemma 4.8.7. Let $m=\operatorname{dim} M$ and let $k=\operatorname{dim} V$. Then there exist universal constants $c_{i}(m, k)$ so that if $P$ is as in Lemma 4.8.1,
(a) $a_{2}(x, P)=(4 \pi)^{-m / 2} \operatorname{Tr}\left(c_{1}(m, k) R_{i j i j}+c_{2}(m, k) E\right)$,
(b)

$$
\begin{gathered}
a_{4}(x, P)=(4 \pi)^{-m / 2} \operatorname{Tr}\left(c_{3}(m, k) R_{i j i j ; k k}+c_{4}(m, k) R_{i j i j} R_{k l k l}\right. \\
+c_{5}(m, k) R_{i j i k} R_{l j l k}+c_{6}(m, k) R_{i j k l} R_{i j k l} \\
+c_{7}(m, k) E_{; i i}+c_{8}(m, k) E^{2} \\
\left.+c_{9}(m, k) E R_{i j i j}+c_{10}(m, k) \Omega_{i j} \Omega_{i j}\right) .
\end{gathered}
$$

This is an important simplification because it shows in particular that terms such as $\operatorname{Tr}(E)^{2}$ do not appear in $a_{4}$.

The first observation we shall need is the following:
Lemma 4.8.8. The constants $c_{i}(m, k)$ of Lemma 4.8 .7 can be chosen to be independent of the dimension $m$ and the fiber dimension $k$.

Proof: The leading symbol of $P$ is scalar. The analysis of Chapter 1 in this case immediately leads to a combinatorial formula for the coefficients in terms of certain trignometric integrals $\int \xi^{\alpha} e^{-|\xi|^{2}} d \xi$ and the fiber dimension does not enter. Alternatively, we could use the additivity of $e_{n}\left(x, P_{1} \oplus P_{2}\right)=$ $e_{n}\left(x, P_{1}\right) \oplus e_{n}\left(x, P_{2}\right)$ of Lemma 1.7.5 to conclude the formulas involved must be independent of the dimension $k$. We may therefore write $c_{i}(m, k)=$ $c_{i}(m)$.

There is a natural restriction map $r: \mathcal{S}_{m, n, k} \rightarrow \mathcal{S}_{m-1, n, k}$ defined by restricting to operators of the form $P=P_{1} \otimes 1+I_{V} \otimes\left(-\partial^{2} / \partial \theta^{2}\right)$ over $M=M_{1} \times S^{1}$. Algebraically, we simply set to zero any variables involving the last index. The multiplicative property of Lemma 1.7.5 implies

$$
r\left(S_{m, n, k}\right)=\sum_{p+q=n} S_{m-1, p, k}\left(P_{1}\right) \otimes S_{1, q, 1}\left(-\partial^{2} / \partial \theta^{2}\right) .
$$

Since all the jets of the symbol of $-\partial^{2} / \partial \theta^{2}$ vanish for $q>0, S_{1, q, 1}=0$ for $q>0$ and $a_{0}=(4 \pi)^{-1 / 2}$ by Lemma 4.8.4. Therefore

$$
r\left(S_{m, n, k}\right)=(4 \pi)^{-1 / 2} S_{m-1, n, k}
$$

Since we have included the normalizing constant $(4 \pi)^{-m / 2}$ in our definition, the constants are independent of the dimension $m$ for $m \geq 4$. If $m=1,2,3$, then the invariants of Lemma 4.8.6 are not linearly independent so we choose the constants to agree with $c_{i}(m, k)$ in these cases.

We remark that if $P$ is a higher order operator with leading symbol given by a power of the metric tensor, then there a similar theory expressing $a_{n}$ in terms of invariant tensorial expressions. However, in this case, the coefficients depend upon the dimension $m$ in a much more fundamental way than simply $(4 \pi)^{-m / 2}$ and we refer to (Gilkey, the spectral geometry of the higher order Laplacian) for further details.

Since the coefficients do not depend on $(m, k)$, we drop the somewhat cumbersome notation $S_{m, n, k}$ and return to the notation $e_{n}(x, P)$ discussed in the first chapter so $\operatorname{Tr}\left(e_{n}(x, P)\right)=a_{n}(x, P)$. We use the properties of the exponential function to compute:

Lemma 4.8.9. Using the notation of Lemma 4.8.7,

$$
c_{2}=1, \quad c_{8}=\frac{1}{2}, \quad c_{9}=c_{1}
$$

Proof: Let $P$ be as in Lemma 4.8.1 and let $a$ be a real constant. We construct the operator $P_{a}=P-a$. The metric and connection are unchanged; we must replace $E$ by $E+a$. Since $e^{-t(P-a)}=e^{-t P} e^{t a}$ we conclude:

$$
e_{n}(x, P-a)=\sum_{p+q=n} e_{p}(x, P) a^{q} / q!
$$

by comparing terms in the asymptotic expansion. We shall ignore factors of $(4 \pi)^{-m / 2}$ henceforth for notational convenience. Then:

$$
e_{2}(x, P-a)=e_{2}(x, P)+c_{2} a=e_{2}(x, P)+e_{0}(x, P) a=e_{2}(x, P)+a
$$

This implies that $c_{2}=1$ as claimed. Next, we have

$$
\begin{aligned}
e_{4}(x, P-a) & =e_{4}(x, P)+c_{8} a^{2}+2 c_{8} a E+c_{9} a R_{i j i j} \\
& =e_{4}(x, P)+e_{2}(x, P) a+e_{0}(x, P) a^{2} / 2 \\
& =e_{4}(x, P)+\left(c_{1} R_{i j i j}+c_{2} E\right) a+a^{2} / 2
\end{aligned}
$$

which implies $c_{8}=1 / 2$ and $c_{1}=c_{9}$ as claimed.
We now use some recursion relations derived in (Gilkey, Recursion relations and the asymptotic behavior of the eigenvalues of the Laplacian). To illustrate these, we first suppose $m=1$. We consider two operators:

$$
A=\partial / \partial x+b, \quad A^{*}=-\partial / \partial x+b
$$

where $b$ is a real scalar function. This gives rise to operators:

$$
\begin{aligned}
& P_{1}=A^{*} A=-\left(\partial^{2} / \partial x^{2}+\left(b^{\prime}-b^{2}\right)\right) \\
& P_{2}=A A^{*}=-\left(\partial^{2} / \partial x^{2}+\left(-b^{\prime}-b^{2}\right)\right)
\end{aligned}
$$

acting on $C^{\infty}\left(S^{1}\right)$. The metric and connection defined by these operators is flat. $E\left(P_{1}\right)=b^{\prime}-b^{2}$ and $E\left(P_{2}\right)=-b^{\prime}-b^{2}$.

Lemma 4.8.10.

$$
(n-1)\left(e_{n}\left(x, P_{1}\right)-e_{n}\left(x, P_{2}\right)\right)=\partial / \partial x\{\partial / \partial x+2 b\} e_{n-2}\left(x, P_{1}\right)
$$

Proof: Let $\left\{\lambda_{\nu}, \theta_{\nu}\right\}$ be a complete spectral resolution of $P_{1}$. We ignore any possible zero spectrum since it won't contribute to the series we shall
be constructing. Then $\left\{\lambda_{\nu}, A \theta_{\nu} / \sqrt{\lambda_{\nu}}\right\}$ is a complete spectral resolution of $P_{2}$. We compute:

$$
\begin{aligned}
\frac{d}{d t}\{K(t, x, x, & \left.\left.P_{1}\right)-K\left(t, x, x, P_{2}\right)\right\} \\
& =\sum e^{-t \lambda_{\nu}}\left\{-\lambda_{\nu} \theta_{\nu} \theta_{\nu}+A \theta_{\nu} A \theta_{\nu}\right\} \\
& =\sum e^{-t \lambda_{\nu}}\left\{-P_{1} \theta_{\nu} \cdot \theta_{\nu}+A \theta_{\nu} A \theta_{\nu}\right\} \\
& =\sum e^{-t \lambda_{\nu}}\left\{\theta_{\nu}^{\prime \prime} \theta_{\nu}+\left(b^{\prime}-b^{2}\right) \theta_{\nu}^{2}+\theta_{\nu}^{\prime} \theta_{\nu}^{\prime}+2 b \theta_{\nu}^{\prime} \theta_{\nu}+b^{2} \theta_{\nu}^{2}\right\} \\
& =\sum e^{-t \lambda_{\nu}}\left(\frac{1}{2} \partial / \partial x\right)(\partial / \partial x+2 b)\left\{\theta_{\nu}^{2}\right\} \\
& =\frac{1}{2} \partial / \partial x(\partial / \partial x+2 b) K\left(t, x, x, P_{1}\right) .
\end{aligned}
$$

We equate terms in the asymptotic expansions

$$
\begin{aligned}
\sum t^{(n-3) / 2} \frac{n-1}{2}\left(e_{n}\left(x, P_{1}\right)\right. & \left.-e_{n}\left(x, P_{2}\right)\right) \\
& \sim \frac{1}{2} \sum t^{(n-1) / 2} \partial / \partial x(\partial / \partial x+2 b) e_{n}\left(x, P_{1}\right)
\end{aligned}
$$

to complete the proof of the lemma.
We apply this lemma to compute the coefficient $c_{7}$. If $n=4$, then we conclude:

$$
\begin{aligned}
& e_{4}\left(x, P_{1}\right)=c_{7}\left(b^{\prime}-b^{2}\right)^{\prime \prime}+c_{8}\left(b^{\prime}-b^{2}\right)^{2}=c_{7} b^{\prime \prime \prime}+\text { lower order terms } \\
& e_{4}\left(x, P_{2}\right)=c_{7}\left(-b^{\prime}-b^{2}\right)^{\prime \prime}+c_{8}\left(-b^{\prime}-b^{2}\right)^{2}=-c_{7} b^{\prime \prime \prime}+\text { lower order terms } \\
& e_{2}\left(x, P_{1}\right)=b^{\prime}-b^{2}
\end{aligned}
$$

so that:

$$
\begin{aligned}
3\left(e_{4}\left(x, P_{1}\right)-e_{4}\left(x, P_{2}\right)\right) & =6 c_{7} b^{\prime \prime \prime}+\text { lower order terms } \\
\partial / \partial x(\partial / \partial x+2 b)\left(b^{\prime}-b^{2}\right) & =b^{\prime \prime \prime}+\text { lower order terms }
\end{aligned}
$$

from which it follows that $c_{7}=1 / 6$. It is also convenient at this stage to obtain information about $e_{6}$. If we let $e_{6}=c E^{\prime \prime \prime \prime}+$ lower order terms then we express:

$$
\begin{aligned}
e_{6}\left(x, P_{1}\right)-e_{6}\left(x, P_{2}\right) & =2 c b^{(5)} \\
\partial / \partial x(\partial / \partial x+2 b)\left(e_{4}\right) & =c_{7} b^{(5)}
\end{aligned}
$$

from which it follows that the constant $c$ is $(4 \pi)^{-1 / 2} \cdot c_{7} / 10=(4 \pi)^{-1 / 2} / 60$. We summarize these results as follows:

Lemma 4.8.11. We can expand $a_{n}$ in the form:
(a) $a_{2}(x, P)=(4 \pi)^{-m / 2} \operatorname{Tr}\left(c_{1} R_{i j i j}+E\right)$.
(b)

$$
\begin{gathered}
a_{4}(x, P)=(4 \pi)^{-m / 2} \operatorname{Tr}\left(c_{3} R_{i j i j ; k k}+c_{4} R_{i j i j} R_{k l k l}+c_{5} R_{i j i k} R_{l j l k}\right. \\
+c_{6} R_{i j k l} R_{i j k l}+E_{; k k} / 6+E^{2} / 2 \\
\left.+c_{1} R_{i j i j} E+c_{10} \Omega_{i j} \Omega_{i j}\right)
\end{gathered}
$$

(c)

$$
\begin{aligned}
a_{6}(x, P)=(4 \pi)^{-m / 2} \operatorname{Tr}\left(E_{; k k l l}\right. & / 60 \\
& \left.+c_{11} R_{i j i j ; k k l l}+\text { lower order terms }\right)
\end{aligned}
$$

Proof: (a) and (b) follow immediately from the computations previously performed. To prove (c) we argue as in the proof of 4.8.7 to show $a_{6}(x, P)=(4 \pi)^{-m / 2} \operatorname{Tr}\left(c E_{; k k}+c_{11} R_{i j i j ; k k}+\right.$ lower order terms $)$ and then use the evaluation of $c$ given above.

We can use a similar recursion relation if $m=2$ to obtain further information regarding these coefficients. We consider the de Rham complex, then:

Lemma 4.8.12. If $m=2$ and $\Delta_{p}$ is the Laplacian on $C^{\infty}\left(\Lambda^{p}\left(T^{*} M\right)\right)$, then:

$$
\frac{(n-2)}{2}\left\{a_{n}\left(x, \Delta_{0}\right)-a_{n}\left(x, \Delta_{1}\right)+a_{n}\left(x, \Delta_{2}\right)\right\}=a_{n-2}\left(x, \Delta_{0}\right)_{; k k}
$$

Proof: This recursion relationship is due to McKean and Singer. Since the invariants are local, we may assume $M$ is orientable and $a_{n}\left(x, \Delta_{0}\right)=$ $a_{n}\left(x, \Delta_{2}\right)$. Let $\left\{\lambda_{\nu}, \theta_{\nu}\right\}$ be a spectral resolution for the non-zero spectrum of $\Delta_{0}$ then $\left\{\lambda_{\nu}, * \theta_{\nu}\right\}$ is a spectral resolution for the non-zero spectrum of $\Delta_{2}$ and $\left\{\lambda_{\nu}, d \theta_{\nu} / \sqrt{\lambda_{\nu}}, \delta * \theta_{\nu} / \sqrt{\lambda_{\nu}}\right\}$ is a spectral resolution for the non-zero spectrum of $\Delta_{1}$. Therefore:

$$
\begin{aligned}
\frac{d}{d t}(K(t, x, x & \left.\left., \Delta_{0}\right)-K\left(t, x, x, \Delta_{1}\right)+K\left(t, x, x, \Delta_{2}\right)\right) \\
& =\sum e^{-t \lambda_{\nu}}\left(-2 \lambda_{\nu} \theta_{\nu} \theta_{\nu}+d \theta_{\nu} \cdot d \theta_{\nu}+\delta * \theta_{\nu} \cdot \delta * \theta_{\nu}\right) \\
& =\sum e^{-t \lambda_{\nu}}\left(-2 \Delta_{0} \theta_{\nu} \cdot \theta_{\nu}+2 d \theta_{\nu} \cdot d \theta_{\nu}\right)=K\left(t, x, x, \Delta_{0}\right)_{; k k}
\end{aligned}
$$

from which the desired identity follows.
Before using this identity, we must obtain some additional information about $\Delta_{1}$.

Lemma 4.8.13. Let $m$ be arbitrary and let $\rho_{i j}=-R_{i j i k}$ be the Ricci tensor. Then $\Delta \theta=-\theta_{; k k}+\rho(\theta)$ for $\theta \in C^{\infty}\left(\Lambda^{1}\right)$ and where $\rho(\theta)_{i}=\rho_{i j} \theta_{j}$. Thus $E\left(\Delta_{2}\right)=-\rho$.

Proof: We apply Lemma 4.8.3 to conclude

$$
E(\theta)=\frac{1}{2} R_{a b i j} \theta_{i} e_{b} * e_{a} * e_{j}
$$

from which the desired result follows using the Bianchi identities.
We can now check at least some of these formulas. If $m=2$, then $E\left(\Delta_{1}\right)=R_{1212} I$ and $E\left(\Delta_{0}\right)=E\left(\Delta_{2}\right)=0$. Therefore:

$$
\begin{aligned}
a_{2}\left(x, \Delta_{0}\right)-a_{2}\left(x, \Delta_{1}\right)+a_{2}\left(x, \Delta_{2}\right) & =(4 \pi)^{-1}\left\{(1-2+1) c_{1} R_{i j i j}-2 R_{1212}\right\} \\
& =-R_{1212} / 2 \pi
\end{aligned}
$$

which is, in fact, the integrand of the Gauss-Bonnet theorem. Next, we compute, supressing the factor of $(4 \pi)^{-1}$ :

$$
\begin{aligned}
a_{4}\left(x, \Delta_{0}\right)-a_{4}\left(x, \Delta_{1}\right)+ & a_{4}\left(x, \Delta_{2}\right) \\
= & \left(c_{3} R_{i j i j ; k k}+c_{4} R_{i j i j} R_{k l k l}+c_{5} R_{i j i k} R_{l j l k}\right. \\
& \left.\quad+c_{6} R_{i j k l} R_{i j k l}\right)(1-2+1) \\
& \quad+\frac{1}{6}(-2)\left(R_{1212 ; k k}\right)-\frac{2}{2} R_{1212} R_{1212} \\
& \quad-4 c_{1} R_{1212} R_{1212}-c_{10} \operatorname{Tr}\left(\Omega_{i j} \Omega_{i j}\right) \\
= & -\frac{1}{3}\left(R_{1212 ; k k}\right)-\left(1+4 c_{1}-4 c_{10}\right)\left(R_{1212}\right)^{2}
\end{aligned}
$$

so that Lemma 4.2.11 applied to the case $n=4$ implies:

$$
a_{4}\left(x, \Delta_{0}\right)-a_{4}\left(x, \Delta_{1}\right)+a_{4}\left(x, \Delta_{2}\right)=a_{2}\left(x, \Delta_{0}\right)_{; k k}
$$

from which we derive the identities:

$$
c_{1}=-\frac{1}{6} \quad \text { and } \quad 1+4 c_{1}-4 c_{10}=0
$$

from which it follows that $c_{10}=1 / 12$. We also consider $a_{6}$ and Lemma 4.8.11(c)

$$
\begin{aligned}
a_{6}\left(x, \Delta_{0}\right)-a_{6}\left(x, \Delta_{1}\right)+a_{6}\left(x, \Delta_{2}\right)= & c_{11} \\
& R_{i j i j ; k k l l}(1-2+1)-\frac{2}{60} R_{1212 ; k k l l} \\
& + \text { lower order terms }
\end{aligned}
$$

so Lemma 4.2.12 implies that $2(-2 / 60)=2 c_{3}$ so that $c_{3}=-1 / 30$.
This leaves only the constants $c_{4}, c_{5}$, and $c_{6}$ undetermined. We let $M=M_{1} \times M_{2}$ be the product manifold with $\Delta_{0}=\Delta_{0}^{1}+\Delta_{0}^{2}$. Then the product formulas of Lemma 1.7.5 imply:

$$
a_{4}\left(\Delta_{0}, x\right)=a_{4}\left(\Delta_{0}^{1}, x_{1}\right)+a_{4}\left(\Delta_{0}^{2}, x_{2}\right)+a_{2}\left(\Delta_{0}^{1}, x_{1}\right) a_{2}\left(\Delta_{0}^{2}, x_{2}\right)
$$

The only term in the expression for $a_{4}$ giving rise to cross terms involving derivatives of both metrics is $c_{4} R_{i j i j} R_{k l k l}$. Consequently, $2 c_{4}=c_{1}^{2}=1 / 36$ so $c_{4}=1 / 72$. We summarize these computations as follows:
Lemma 4.8.14. We can expand $\alpha_{n}$ in the form:
(a) $a_{2}(x, P)=(4 \pi)^{-m / 2} \operatorname{Tr}\left(-R_{k j k j}+6 E\right) / 6$.
(b)

$$
\begin{aligned}
a_{4}(x, P)= & \frac{(4 \pi)^{-m / 2}}{360} \operatorname{Tr}\left(-12 R_{i j i j ; k k}+5 R_{i j i j} R_{k l k l}+c_{5} R_{i j i k} R_{l j l k}\right. \\
& \left.+c_{6} R_{i j k l} R_{i j k l}-60 E R_{i j i j}+180 E^{2}+60 E_{; k k}+30 \Omega_{i j} \Omega_{i j}\right)
\end{aligned}
$$

We have changed the notation slightly to introduce the common denominator 360 .

We must compute the universal constants $c_{5}$ and $c_{6}$ to complete our determination of the formula $a_{4}$. We generalize the recursion relations of Lemmas 4.8.10 and 4.8.12 to arbitrary dimensions as follows. Let $M=T_{m}$ be the $m$-dimensional torus with usual periodic parameters $0 \leq x_{i} \leq 2$ for $i=1, \ldots, m$. We let $\left\{e_{i}\right\}$ be a collection of Clifford matrices so $e_{i} e_{j}+e_{j} e_{i}=$ $2 \delta_{i j}$. Let $h(x)$ be a real-valued function on $T_{m}$ and let the metric be given by:

$$
d s^{2}=e^{-h}\left(d x_{1}^{2}+\cdots+d x_{m}^{2}\right), \quad \mathrm{dvol}=e^{-h m / 2} d x_{1} \ldots d x_{m}
$$

We let the operator $A$ and $A^{*}$ be defined by:

$$
\begin{aligned}
A & =e^{m h / 4} \sum e_{j} \frac{\partial}{\partial x_{j}} e^{(2-m) h / 4} \\
A^{*} & =e^{(2+m) h / 4} \sum e_{j} \frac{\partial}{\partial x_{j}} e^{-m h / 4}
\end{aligned}
$$

and define:

$$
\begin{aligned}
& P_{1}=A^{*} A=-e^{(2+m) h / 4} \sum_{j} \frac{\partial^{2}}{\partial x_{j}^{2}} e^{(2-m) h / 4} \\
&=-e^{h}\left\{\sum \frac{\partial^{2}}{\partial x_{j}^{2}}+\frac{1}{2}(2-m) h_{/ j} \frac{\partial}{\partial x_{j}}\right. \\
&\left.+\frac{1}{16}\left(4(2-m) h_{/ j j}+(2-m)^{2} h_{/ j} h_{/ j}\right)\right\} \\
& P_{2}=A A^{*}= e^{m h / 4} \sum_{j} e_{j} \frac{\partial}{\partial x_{j}} e^{h} \sum_{k} e_{k} \frac{\partial}{\partial x_{k}} e^{-m h / 4}
\end{aligned}
$$

Lemma 4.8.15. With the notation given above,

$$
(n-m)\left\{a_{n}\left(x, P_{1}\right)-a_{n}\left(x, P_{2}\right)\right\}=e^{h m / 2} \frac{\partial^{2}}{\partial x_{k}^{2}} e^{h(2-m) / 2} a_{n-2}(x, P)
$$

Proof: We let $P_{0}$ be the scalar operator

$$
-e^{(2+m) h / 4} \sum_{j} \frac{\partial^{2}}{\partial x_{j}^{2}} e^{(2-m) h / 4}
$$

If the representation space on which the $e_{j}$ act has dimension $u$, then $P_{1}=P_{0} \otimes I_{u}$. We let $v_{s}$ be a basis for this representation space. Let $\left\{\lambda_{\nu}, \theta_{\nu}\right\}$ be a spectral resolution for the operator $P_{0}$ then $\left\{\lambda_{\nu}, \theta_{\nu} \otimes v_{s}\right\}$ is a spectral resolution of $P_{1}$. We compute:

$$
\begin{aligned}
& \frac{d}{d t}\left(\operatorname{Tr} K\left(t, x, x, P_{1}\right)-\operatorname{Tr}\left(t, x, x, P_{2}\right)\right) \\
& \quad=\sum_{\nu, s} e^{-t \lambda_{\nu}}\left\{-\left(P_{1} \theta_{\nu} \otimes u_{s}, \theta_{\nu} \otimes u_{s}\right)+\left(A \theta_{\nu} \otimes u_{s}, A \theta_{\nu} \otimes u_{s}\right)\right\}
\end{aligned}
$$

where (, ) denotes the natural inner product $\left(u_{s}, u_{s^{\prime}}\right)=\delta_{s, s^{\prime}}$. The $e_{j}$ are self-adjoint matrices. We use the identity $\left(e_{j} u_{s}, e_{k} u_{s}\right)=\delta_{j k}$ to compute:

$$
\begin{aligned}
& =\sum_{\nu} v \cdot e^{-t \lambda_{\nu}} e^{h}\left\{\theta_{\nu / k k} \theta_{\nu}+\frac{1}{2}(2-m) h_{/ k} \theta_{\nu / k} \theta_{\nu}+\frac{1}{4}(2-m) h_{/ k k} \theta_{\nu} \theta_{\nu}\right. \\
& \quad+\frac{1}{16}(2-m)^{2} h_{/ k} h_{/ k} \theta_{\nu} \theta_{\nu}+\theta_{\nu / k} \theta_{\nu / k}+\frac{1}{2}(2-m) h_{/ k} \theta_{\nu / k} \theta_{\nu} \\
& \\
& \left.\quad+\frac{1}{16}(2-m)^{2} h_{/ k} h_{/ k} \theta_{\nu} \theta_{\nu}\right\} \\
& =\sum_{\nu} v \cdot e^{-t \lambda_{\nu}} e^{h m / 2} \frac{1}{2} \sum_{j} \frac{\partial^{2}}{\partial x_{j}^{2}} e^{h(2-m) / 2}\left(\theta_{\nu} \theta_{\nu}\right) \\
& = \\
& \frac{1}{2} e^{h m / 2} \sum_{j} \frac{\partial^{2}}{\partial x_{j}^{2}} e^{h(2-m) / 2} \operatorname{Tr} K\left(t, x, x, P_{1}\right)
\end{aligned}
$$

We compare coefficients of $t$ in the two asymptotic expansions to complete the proof of the lemma.

We apply Lemma 4.8 .15 to the special case $n=m=6$. This implies that

$$
\sum_{j} \frac{\partial^{2}}{\partial x_{j}^{2}} e^{-2 h} a_{4}\left(x, P_{1}\right)=0
$$

Since $a_{x}\left(x, P_{1}\right)$ is a formal polynomial in the jets of $h$ with coefficients which are smooth functions of $h$, and since $a_{4}$ is homogeneous of order 4,
this identity for all $h$ implies $a_{4}\left(x, P_{1}\right)=0$. This implies $a_{4}\left(x, P_{0}\right)=$ $\frac{1}{v} a_{4}\left(x, P_{1}\right)=0$.

Using the formulas of lemmas 4.8.1 and 4.8.2 it is an easy exercise to calculate $R_{i j k l}=0$ if all 4 indices are distinct. The non-zero curvatures are given by:

$$
\begin{array}{rlrl}
R_{i j i}^{k} & =-\frac{1}{4}\left(2 h_{/ j k}+h_{/ j} h_{/ k}\right) & \text { if } j \neq k & \\
R_{i j i}^{j} & =-\frac{1}{4}\left(2 h_{/ i i}+2 h_{/ j j}-\sum_{k \neq i, k \neq j} h_{/ k} h_{/ k}\right. & & (\text { don't sum over } i) \\
\Omega_{i j} & =0 & & \\
E\left(P_{0}\right) & =e^{h} \sum_{k}\left(-h_{/ k k}+h_{/ k} h_{/ k}\right) . &
\end{array}
$$

When we contract the curvature tensor to form scalar invariants, we must include the metric tensor since it is not diagonal. This implies:

$$
\tau=\sum_{i, j} g^{i i} R_{i j i}^{j}=e^{h} \sum_{k}\left(-5 h_{/ k k}+5 h_{/ k} h_{/ k}\right)=5 E\left(P_{1}\right)
$$

which implies the helpful identities:

$$
-12 \tau_{; k k}+60 E_{; k k}=0 \quad \text { and } \quad 5 R_{i j i j} R_{k l k l}-60 R_{i j i j} E+180 E^{2}=5 E^{2}
$$

so that we conclude:

$$
5 E^{2}+c_{5} R_{i j i k} R_{l j l k}+c_{6} R_{i j k l} R_{i j k l}=0
$$

We expand:

$$
\begin{aligned}
R_{i j i k} R_{l j l k} & =e^{2 h}\left(\frac{15}{2} h_{/ 11} h_{/ 11}+8 h_{/ 12} h_{/ 12}+\text { other terms }\right) \\
R_{i j k l} R_{i j k l} & =e^{2 h}\left(5 h_{/ 11} h_{/ 11}+8 h_{/ 12} h_{/ 12}+\text { other terms }\right) \\
E^{2} & =e^{2 h}\left(h_{/ 11} h_{/ 11}+0 \cdot h_{/ 12} h_{/ 12}+\text { other terms }\right)
\end{aligned}
$$

so we conclude finally:

$$
15 c_{5}+10 c_{6}+10=0 \quad \text { and } \quad c_{5}+c_{6}=0
$$

We solve these equations to conclude $c_{5}=-2$ and $c_{6}=2$ which proves finally:

Theorem 4.8.16. Let $P$ be a second order differential operator with leading symbol given by the metric tensor. Let $P=P_{\nabla}-E$ be decomposed as in 4.8.1. Let $a_{n}(x, P)$ be the invariants of the heat equation discussed in Chapter 1.
(a) $a_{0}(x, P)=(4 \pi)^{-m / 2} \operatorname{Tr}(I)$.
(b) $a_{2}(x, P)=(4 \pi)^{-m / 2} \operatorname{Tr}\left(-R_{i j i j}+6 E\right) / 6$.
(c)

$$
\begin{aligned}
& a_{4}(x, P)=\frac{(4 \pi)^{-m / 2}}{360} \times \\
& \begin{aligned}
& \operatorname{Tr}\left(-12 R_{i j i j ; k k}+5 R_{i j i j} R_{k l k l}-2 R_{i j i k} R_{l j l k}+2 R_{i j k l} R_{i j k l}\right. \\
&\left.-60 R_{i j i j} E+180 E^{2}+60 E_{; k k}+30 \Omega_{i j} \Omega_{i j}\right)
\end{aligned}
\end{aligned}
$$

(d)

$$
\begin{aligned}
& a_{6}(x, P)=(4 \pi)^{-m / 2} \times \\
& \begin{array}{c}
\operatorname{Tr}\left\{\begin{array}{r}
\frac{1}{7!}\left(-18 R_{i j i j ; k k l l}+17 R_{i j i j ; k} R_{u l u l ; k}-2 R_{i j i k ; l} R_{u j u k ; l}\right. \\
-4 R_{i j i k ; l} R_{u j u l ; k}+9 R_{i j k u ; l} R_{i j k u ; l}+28 R_{i j i j} R_{k u k u ; l l} \\
\left.\quad-8 R_{i j i k} R_{u j u k ; l l}+24 R_{i j i k} R_{u j u l ; k l}+12 R_{i j k l} R_{i j k l ; u u}\right)
\end{array}\right. \\
+\frac{1}{9 \cdot 7!}\left(-35 R_{i j i j} R_{k l k l} R_{p q p q}+42 R_{i j i j} R_{k l k p} R_{q l q p}\right. \\
\quad-42 R_{i j i j} R_{k l p q} R_{k l p q}+208 R_{i j i k} R_{j u l u} R_{k p l p} \\
\quad-192 R_{i j i k} R_{u p l p} R_{j u k l}+48 R_{i j i k} R_{j u l p} R_{k u l p} \\
\left.\quad-44 R_{i j k u} R_{i j l p} R_{k u l p}-80 R_{i j k u} R_{i l k p} R_{j l u p}\right) \\
+\frac{1}{360}\left(8 \Omega_{i j ; k} \Omega_{i j ; k}+2 \Omega_{i j ; j} \Omega_{i k ; k}+12 \Omega_{i j} \Omega_{i j ; k k}-12 \Omega_{i j} \Omega_{j k} \Omega_{k i}\right. \\
\left.-6 R_{i j k l} \Omega_{i j} \Omega_{k l}+4 R_{i j i k} \Omega_{j l} \Omega_{k l}-5 R_{i j i j} \Omega_{k l} \Omega_{k l}\right) \\
+\frac{1}{360}\left(6 E_{; i i j j}+60 E E_{; i i}+30 E_{; i} E_{; i}+60 E^{3}+30 E \Omega_{i j} \Omega_{i j}\right. \\
\quad-10 R_{i j i j} E_{; k k}-4 R_{i j i k} E_{; j k}-12 R_{i j i j ; k} E_{; k}-30 R_{i j i j} E^{2} \\
\quad-12 R_{i j i j ; k k} E+5 R_{i j i j} R_{k l k l} E \\
\left.\left.\quad-2 R_{i j i k} R_{i j k l} E+2 R_{i j k l} R_{i j k l} E\right) .\right\}
\end{array}
\end{aligned}
$$

Proof: We have derived (a)-(c) explicitly. We refer to (Gilkey, The spectral geometry of a Riemannian manifold) for the proof of (d) as it is quite long and complicated. We remark that our sign convention is that $R_{1212}=-1$ on the sphere of radius 1 in $R^{3}$.

We now begin our computation of $a_{n}\left(x, \Delta_{p}^{m}\right)$ for $n=0,2,4$.

Lemma 4.8.17. Let $m=4$, and let $\Delta_{p}$ be the Laplacian on $p$-forms. Decompose

$$
\begin{aligned}
& a_{2}\left(x, \Delta_{p}\right)=(4 \pi)^{-2} \cdot c_{0}(p) \cdot R_{i j i j} / 6 \\
& a_{4}\left(x, \Delta_{p}\right)=(4 \pi)^{-2}\left\{c_{1}(p) R_{i j i j ; k k}+c_{2}(p) R_{i j i j} R_{k l k l}\right. \\
& \left.\qquad c_{3}(p) R_{i j i k} R_{l j l k}+c_{4}(p) R_{i j k l} R_{i j k l}\right\} / 360
\end{aligned}
$$

Then the $c_{i}$ are given by the following table:

|  | $c_{0}$ | $c_{1}$ | $c_{2}$ | $c_{3}$ | $c_{4}$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| $p=0$ | -1 | -12 | 5 | -2 | 2 |
| $p=1$ | 2 | 12 | $-40$ | 172 | $-22$ |
| $p=2$ | 6 | 48 | 90 | $-372$ | 132 |
| $p=3$ | 2 | 12 | -40 | 172 | -22 |
| $p=4$ | -1 | -12 | 5 | -2 | 2 |
| $\sum(-1)^{p}$ | 0 | 0 | 180 | $-720$ | 180 |

Proof: By Poincare duality, $a_{n}\left(x, \Delta_{p}\right)=a_{n}\left(x, \Delta_{4-p}\right)$ so we need only check the first three rows as the corresponding formulas for $\Delta_{3}$ and $\Delta_{4}$ will follow. In dimension 4 the formula for the Euler form is $(4 \pi)^{-2}\left(R_{i j i j} R_{k l k l}-\right.$ $\left.4 R_{i j i k} R_{l j l k}+R_{i j k l} R_{i j k l}\right) / 2$ so that the last line follows from Theorem 2.4.8.

If $p=0$, then $E=\Omega=0$ so the first line follows from Theorem 4.8.15. If $p=1$, then $E=-\rho_{i j}=R_{i k i j}$ is the Ricci tensor by Lemma 4.8.13. Therefore:

$$
\begin{aligned}
\operatorname{Tr}\left(-12 R_{i j i j ; k k}+60 E_{; k k}\right) & =-48 R_{i j i j ; k k}+60 R_{i j i j ; k k} \\
& =12 R_{i j i j ; k k} \\
\operatorname{Tr}\left(5 R_{i j i j} R_{k l k l}-60 R_{i j i j} E\right) & =20 R_{i j i j} R_{k l k l}-60 R_{i j i j} R_{k l k l} \\
& =-40 R_{i j i j} R_{k l k l} \\
\operatorname{Tr}\left(-2 R_{i j i k} R_{l k l k}+180 E^{2}\right) & =-8 R_{i j i k} R_{l j l k}+180 R_{i j i k} R_{l j l k} \\
& =172 R_{i j i k} R_{l j l k} \\
\operatorname{Tr}\left(2 R_{i j k l} R_{i j k l}+30 \Omega_{i j} \Omega_{i j}\right) & =8 R_{i j k l} R_{i j k l}-30 R_{i j k l} R_{i j k l} \\
& =-22 R_{i j k l} R_{i j k l}
\end{aligned}
$$

which completes the proof of the second line. Thus the only unknown is $c_{k}(2)$. This is computed from the alternating sum and completes the proof.

More generally, we let $m>4$. Let $M=M_{4} \times T_{m-4}$ be a product manifold, then this defines a restriction map $r_{m-4}: \mathcal{P}_{m, n} \rightarrow \mathcal{P}_{4, n}$ which is
an isomorphism for $n \leq 4$. Using the multiplication properties given in Lemma 1.7.5, it follows:

$$
a_{n}\left(x, \Delta_{p}^{m}\right)=\sum_{\substack{i+j=n \\ p_{1}+p_{2}=p}} a_{i}\left(x, \Delta_{p_{1}}^{4}\right) a_{j}\left(x, \Delta_{p_{2}}^{m-4}\right)
$$

On the flat torus, all the invariants vanish for $j>0$. By Lemma 4.8.4 $a_{0}\left(x, \Delta_{p_{2}}^{m-4}\right)=(4 \pi)^{-m / 2}\binom{m-4}{p_{2}}$ so that:

$$
a_{n}\left(x, \Delta_{u}^{4}\right)=\sum_{u+v=p}(4 \pi)^{(m-4) / 2}\binom{m-4}{v} a_{n}\left(x, \Delta_{u}^{4}\right) \quad \text { for } n \leq 4
$$

where $a_{n}\left(x, \Delta_{u}^{4}\right)$ is given by Lemma 4.8.16. If we expand this as a polynomial in $m$ for small values of $u$ we conclude:

$$
\begin{aligned}
a_{2}\left(x, \Delta_{1}^{m}\right)= & \frac{(4 \pi)^{-m / 2}}{6}(6-m) R_{i j i j} \\
a_{4}\left(x, \Delta_{1}^{m}\right)= & \frac{(4 \pi)^{-m / 2}}{360}\left\{(60-12 m) R_{i j i j ; k k}+(5 m-60) R_{i j i j} R_{k l k l}\right. \\
& \left.\quad+(180-2 m) R_{i j i k} R_{l j l k}+(2 m-30) R_{i j k l} R_{i j k l}\right\}
\end{aligned}
$$

and similarly for $a_{2}\left(x, \Delta_{2}^{m}\right)$ and $a_{4}\left(x, \Delta_{2}^{m}\right)$. In this form, the formulas also hold true for $m=2,3$. We summarize our conclusions in the following theorem:

ThEOREM 4.8.18. Let $\Delta_{p}^{m}$ denote the Laplacian acting on the space of smooth p-forms on an m-dimensional manifold. We let $R_{i j k l}$ denote the curvature tensor with the sign convention that $R_{1212}=-1$ on the sphere of radius 1 in $\mathbf{R}^{3}$. Then:
(a) $a_{0}\left(x, \Delta_{p}^{m}\right)=(4 \pi)^{-m / 2}\binom{m}{p}$.
(b) $a_{2}\left(x, \Delta_{p}^{m}\right)=\frac{(4 \pi)^{-m / 2}}{6}\left\{\binom{m-2}{p-2}+\binom{m-2}{p}-4\binom{m-2}{p-1}\right\}\left(-R_{i j i j}\right)$.
(c) Let

$$
\begin{aligned}
& a_{4}\left(x, \Delta_{p}^{m}\right)=\frac{(4 \pi)^{-m / 2}}{360}\left(c_{1}(m, p) R_{i j i j ; k k}+c_{2}(m, p) R_{i j i j} R_{k l k l}\right. \\
& \left.\quad+c_{3}(m, p) R_{i j i k} R_{l j l k}+c_{4}(m, p) R_{i j k l} R_{i j k l}\right)
\end{aligned}
$$

Then for $m \geq 4$ the coefficients are:

$$
\begin{aligned}
& c_{1}(m, p)=-12\left[\binom{m-4}{p}+\binom{m-4}{p-4}\right]+12\left[\binom{m-4}{p-1}+\binom{m-4}{p-3}\right]+48\binom{m-4}{p-2} \\
& c_{2}(m, p)=5\left[\binom{m-4}{p}+\binom{m-4}{p-4}\right]-40\left[\binom{m-4}{p-1}+\binom{m-4}{p-3}\right]+90\binom{m-4}{p-2} \\
& c_{3}(m, p)=-2\left[\binom{m-4}{p}+\binom{m-4}{p-4}\right]+172\left[\binom{m-4}{p-1}+\binom{m-4}{p-3}\right]-372\binom{m-4}{p-2} \\
& c_{4}(m, p)=2\left[\binom{m-4}{p}+\binom{m-4}{p-4}\right]-22\left[\binom{m-4}{p-1}+\binom{m-4}{p-3}\right]+132\binom{m-4}{p-2} .
\end{aligned}
$$

(d)

$$
\begin{aligned}
a_{0}\left(x, \Delta_{0}^{m}\right) & =(4 \pi)^{-m / 2} \\
a_{2}\left(x, \Delta_{0}^{m}\right)= & \frac{(4 \pi)^{-m / 2}}{6}\left(-R_{i j i j}\right) \\
a_{4}\left(x, \Delta_{0}^{m}\right)= & \frac{(4 \pi)^{-m / 2}}{360}\left(-12 R_{i j i j ; k k}+5 R_{i j i j} R_{k l k l}-2 R_{i j i k} R_{l j l k}\right. \\
& \left.+2 R_{i j k l} R_{i j k l}\right)
\end{aligned}
$$

(e)

$$
\begin{aligned}
a_{0}\left(x, \Delta_{1}^{m}\right)= & (4 \pi)^{-m / 2} \\
a_{2}\left(x, \Delta_{1}^{m}\right)= & \frac{(4 \pi)^{-m / 2}}{6}(6-m) R_{i j i j} \\
a_{4}\left(x, \Delta_{1}^{m}\right)= & \frac{(4 \pi)^{-m / 2}}{360}\left\{(60-12 m) R_{i j i j ; k k}+(5 m-60) R_{i j i j} R_{k l k l}\right. \\
& \quad+(180-2 m) R_{i j i k} R_{l j l k} \\
& \left.\left.\quad+(2 m-30) R_{i j k l} R_{i j k l}\right)\right\}
\end{aligned}
$$

(f)

$$
\begin{aligned}
& a_{0}\left(x, \Delta_{2}^{m}\right)=\frac{(4 \pi)^{-m / 2}}{2} m(m-1) \\
& a_{2}\left(x, \Delta_{2}^{m}\right)=\frac{(4 \pi)^{-m / 2}}{12}\left(-m^{2}+13 m-24\right) R_{i j i j} \\
& a_{4}\left(x, \Delta_{2}^{m}\right)=\frac{(4 \pi)^{-m / 2}}{720}\left\{\left(-12 m^{2}+108 m-144\right) R_{i j i j k ; k k}\right. \\
& +\left(5 m^{2}-115 m+560\right) R_{i j i j} R_{k l k l} \\
& +\left(-2 m^{2}+358 m-2144\right) R_{i j i k} R_{l j l k} \\
& \left.+\left(2 m^{2}-58 m+464\right) R_{i j k l} R_{i j k l}\right\}
\end{aligned}
$$

These results are, of course, not new. They were first derived by Patodi. We could apply similar calculations to determine $a_{0}, a_{2}$, and $a_{4}$ for any operator which is natural in the sense of Epstein and Stredder. In particular, the Dirac operator can be handled in this way.

In principal, we could also use these formulas to compute $a_{6}$, but the lower order terms become extremely complicated. It is not too terribly difficult, however, to use these formulas to compute the terms in $a_{6}$ which involve the 6 jets of the metric and which are bilinear in the 4 and 2 jets of the metric. This would complete the proof of the result concerning $a_{6}$ if $m=4$ discussed in section 4.7.

### 4.9. Spectral Geometry.

Let $M$ be a compact Riemannian manifold without boundary and let $\operatorname{spec}(M, \Delta)$ denote the spectrum of the scalar Laplacian where each eigenvalue is repeated according to its multiplicity. Two manifolds $M$ and $\bar{M}$ are said to be isospectral if $\operatorname{spec}(M, \Delta)=\operatorname{spec}(\bar{M}, \Delta)$. The leading term in the asymptotic expansion of the heat equation is $(4 \pi)^{-m / 2} \cdot \operatorname{vol}(M) \cdot t^{-m / 2}$ so that if $M$ and $\bar{M}$ are isospectral,

$$
\operatorname{dim} M=\operatorname{dim} \bar{M} \quad \text { and } \quad \operatorname{volume}(M)=\operatorname{volume}(\bar{M})
$$

so these two quantitites are spectral invariants. If $P \in \mathcal{P}_{m, n, 0}$ is an invariant polynomial, we define $P(M)=\int_{M} P(G) \mid$ dvol $\mid$. (This depends on the metric in general.) Theorem 4.8.18 then implies that $R_{i j i j}(M)$ is a spectral invasriant since this appears with a non-zero coefficient in the asymptotic expansion of the heat equation.

The scalar Laplacian is not the only natural differential operator to study. (We use the word "natural" in the technical sense of Epstein and Stredder in this context.) Two Riemannian manifolds $M$ and $\bar{M}$ are said to be strongly isospectral if $\operatorname{spec}(M, P)=\operatorname{spec}(\bar{M}, P)$ for all natural operators $P$. Many of the global geometry properties of the manifold are reflected by their spectral geometry. Patodi, for example, proved:

Theorem 4.9.1 (Patodi). Let $\operatorname{spec}\left(M, \Delta_{p}\right)=\operatorname{spec}\left(\bar{M}, \Delta_{p}\right)$ for $p=$ $0,1,2$. Then:
(a)

$$
\begin{array}{ll}
\operatorname{dim} M=\operatorname{dim} \bar{M}, & \operatorname{volume}(M)=\operatorname{volume}(\bar{M}), \\
R_{i j i j}(M)=R_{i j i j}(\bar{M}), & R_{i j i j} R_{k l k l}(M)=R_{i j i j} R_{k l k l}(\bar{M}), \\
R_{i j i k} R_{l j l k}(M)=R_{i j i k} R_{l j l k}(\bar{M}), & R_{i j k l} R_{i j k l}(M)=R_{i j k l} R_{i j k l}(\bar{M}) .
\end{array}
$$

(b) If $M$ has constant scalar curvature $c$, then so does $\bar{M}$.
(c) If $M$ is Einstein, then so is $\bar{M}$.
(d) If $M$ has constant sectional curvature $c$, then so does $\bar{M}$.

Proof: The first three identities of (a) have already been derived. If $m \geq 4$, the remaining 3 integral invariants are independent. We know:

$$
\begin{gathered}
a_{4}\left(\Delta_{p}^{m}\right)=c_{2}(m, p) R_{i j i j} R_{k l k l}+c_{3}(m, p) R_{i j i k} R_{l j l k}(M) \\
+c_{4}(m, p) R_{i j k l} R_{i j k l}(M)
\end{gathered}
$$

As $p=0,1,2$ the coefficients form a $3 \times 3$ matrix. If we can show the matrix has rank 3 , we can solve for the integral invariants in terms of the
spectral invariants to prove (a). Let $\nu=m-4$ and $c=(4 \pi)^{-\nu / 2}$. then our computations in section 4.8 show

$$
\begin{aligned}
c_{j}(m, 0) & =c \cdot c_{j}(4,0) \\
c_{j}(m, 1) & =c \cdot \nu \cdot c_{j}(4,0)+c \cdot c_{j}(4,1) \\
c_{j}(m, 2) & =c \cdot \nu \cdot(\nu-1) / 2 \cdot c_{j}(4,0)+c \cdot \nu \cdot c_{j}(4,1)+c \cdot c_{j}(4,2)
\end{aligned}
$$

so that the matrix for $m>4$ is obtained from the matrix for $m=4$ by elementary row operations. Thus it suffices to consider the case $m=4$. By 4.8.17, the matrix there (modulo a non-zero normalizing constant) is:

$$
\left(\begin{array}{rrr}
5 & -2 & 2 \\
-40 & 172 & -22 \\
90 & -372 & 132
\end{array}\right)
$$

The determinant of this matrix is non-zero from which (a) follows. The case $m=2$ and $m=3$ can be checked directly using 4.8.18.

To prove (b), we note that $M$ has constant scalar curvature $c$ if and only if $\left(2 c+R_{i j i j}\right)^{2}(M)=0$ which by (a) is a spectral invariant. (c) and (d) are similar.

From (d) follows immediately the corollary:
Corollary 4.9.2. Let $M$ and $\bar{M}$ be strongly isospectral. If $M$ is isometric to the standard sphere of radius $r$, then so is $\bar{M}$.

Proof: If $M$ is a compact manifold with sectional curvature $1 / r$, then the universal cover of $M$ is the sphere of radius $r$. If $\operatorname{vol}(M)$ and $\operatorname{vol}(S(r))$ agree, then $M$ and the sphere are isometric.

There are a number of results which link the spectral and the global geometry of a manifold. We list two of these results below:
Theorem 4.9.3. Let $M$ and $\bar{M}$ be strongly isospectral manifolds. Then:
(a) If $M$ is a local symmetric space (i.e., $\nabla R=0$ ), then so is $\bar{M}$.
(b) If the Ricci tensor of $M$ is parallel (i.e., $\nabla \rho=0$ ), then the Ricci tensor of $\bar{M}$ is parallel. In this instance, the eigenvalues of $\rho$ do not depend upon the particular point of the manifold and they are the same for $M$ and $\bar{M}$.

Although we have chosen to work in the real category, there are also isospectral results available in the holomorphic category:
Theorem 4.9.4. Let $M$ and $\bar{M}$ be holomorphic manifolds and suppose $\operatorname{spec}\left(M, \Delta_{p, q}\right)=\operatorname{spec}\left(\bar{M}, \Delta_{p, q}\right)$ for all $(p, q)$.
(a) If $M$ is Kaehler, then so is $\bar{M}$.
(b) If $M$ is $\mathbf{C} P_{n}$, then so is $\bar{M}$.

At this stage, the natural question to ask is whether or not the spectral geometry completely determines $M$. This question was phrased by Kac in
the form: Can you hear the shape of a drum? It is clear that if there exists an isometry between two manifolds, then they are strongly isospectral. That the converse need not hold was shown by Milnor, who gave examples of isospectral tori which were not isometric. In 1978 Vigneras gave examples of isospectral manifolds of constant negative curvature which are not isometric. (One doesn't know yet if they are strongly isospectral.) If the dimension is at least 3 , then the manifolds have different fundamental groups, so are not homotopic. The fundamental groups in question are all infinite and the calculations involve some fairly deep results in quaternion algebras.

In 1983 Ikeda constructed examples of spherical space forms which were strongly isospectral but not isometric. As de Rham had shown that diffeomorphic spherical space forms are isometric, these examples are not diffeomorphic. Unlike Vigneras' examples, Ikeda's examples involve finite fundamental groups and are rather easily studied. In the remainder of this section, we will present an example in dimension 9 due to Ikeda illustrating this phenomenon. These examples occur much more generally, but this example is particularly easy to study. We refer to (Gilkey, On spherical space forms with meta-cyclic fundamental group which are isospectral but not equivariant cobordant) for more details.

Let $G$ be the group of order 275 generated by two elements $A, B$ subject to the relations:

$$
A^{11}=B^{25}=1 \quad \text { and } \quad B A B^{-1}=A^{3}
$$

We note that $3^{5} \equiv 1 \bmod 11$. This group is the semi-direct product $\mathbf{Z}_{11} \propto \mathbf{Z}_{25}$. The center of $G$ is generated by $B^{5}$ and the subgroup generated by $A$ is a normal subgroup. We have short exact sequences:

$$
0 \rightarrow \mathbf{Z}_{11} \rightarrow G \rightarrow \mathbf{Z}_{25} \rightarrow 0 \quad \text { and } \quad 0 \rightarrow \mathbf{Z}_{55} \rightarrow G \rightarrow \mathbf{Z}_{5} \rightarrow 0
$$

We can obtain an explicit realization of $G$ as a subgroup of $\mathrm{U}(5)$ as follows. Let $\alpha=e^{2 \pi i / 11}$ and $\beta=e^{2 \pi i / 25}$ be primitive roots of unity. Let $\left\{e_{j}\right\}$ be the standard basis for $\mathbf{C}^{5}$ and let $\sigma \in \mathrm{U}(5)$ be the permutation matrix $\sigma\left(e_{j}\right)=e_{j-1}$ where the index $j$ is regarded as defined mod 5 . Define a representation:

$$
\pi_{k}(A)=\operatorname{diag}\left(\alpha, \alpha^{3}, \alpha^{9}, \alpha^{5}, \alpha^{4}\right) \quad \text { and } \quad \pi_{k}(B)=\beta^{k} \cdot \sigma
$$

It is immediate that $\pi_{k}(A)^{11}=\pi_{k}(B)^{25}=1$ and it is an easy computation that $\pi_{k}(B) \pi_{k}(A) \pi_{k}(B)^{-1}=\pi_{k}(A)^{3}$ so this extends to a representation of $G$ for $k=1,2,3,4$. (If $H$ is the subgroup generated by $\left\{A, B^{5}\right\}$, we let $\rho(A)=\alpha$ and $\rho\left(B^{5}\right)=\beta^{5 k}$ be a unitary representation of $H$. The representation $\pi_{k}=\rho^{G}$ is the induced representation.)

In fact these representations are fixed point free:

Lemma 4.9.5. Let $G$ be the group of order 275 generated by $\{A, B\}$ with the relations $A^{11}=B^{25}=1$ and $B A B^{-1}=A^{3}$. The center of $G$ is generated by $B^{5}$; the subgroup generated by $A$ is normal. Let $\alpha=e^{2 \pi i / 11}$ and $\beta=e^{2 \pi i / 25}$ be primitive roots of unity. Define representations of $G$ in $\mathrm{U}(5)$ for $1 \leq k \leq 4$ by: $\pi_{k}(A)=\operatorname{diag}\left(\alpha, \alpha^{3}, \alpha^{9}, \alpha^{5}, \alpha^{4}\right)$ and $\pi_{k}(B)=\beta^{k} \cdot \sigma$ where $\sigma$ is the permutation matrix defined by $\sigma\left(e_{i}\right)=e_{i-1}$.
(a) Enumerate the elements of $G$ in the form $A^{a} B^{b}$ for $0 \leq a<11$ and $0 \leq b<25$. If $(5, b)=1$, then $A^{a} B^{b}$ is conjugate to $B^{b}$.
(b) The representations $\pi_{k}$ are fixed point free.
(c) The eigenvalues of $\pi_{k}\left(A^{a} B^{b}\right)$ and $\pi_{1}\left(A^{a} B^{k b}\right)$ are the same so these two matrices are conjugate in $\mathrm{U}(5)$.

Proof: $A^{-j} B^{b} A^{j}=A^{j\left(3^{b}-1\right)} B^{b}$. If $(5, b)=1$, then $3^{b}-1$ is coprime to 11 so we can solve the congruence $j\left(3^{b}-1\right) \equiv a \bmod 11$ to prove (a). Suppose $(5, b)=1$ so the eigenvalues of $\pi_{k}\left(A^{a} B^{b}\right)$ and $\pi_{k}\left(B^{b}\right)$ coincide. Let $\left\{\varepsilon_{1}, \ldots, \varepsilon_{5}\right\}$ be the $5^{\text {th }}$ roots of unity; these are the eigenvalues of $\sigma$ and of $\sigma^{b}$. Thus the eigenvalues of $\pi_{k}\left(B^{b}\right)$ are $\left\{\beta^{k b} \varepsilon_{1}, \ldots, \beta^{k b} \varepsilon_{5}\right\}$ and are primitive $25^{\text {th }}$ roots of unity. Thus $\operatorname{det}\left(\pi_{k}\left(B^{b}\right)-I\right) \neq 0$ and $\pi_{k}\left(B^{b}\right)$ has the same eigenvalues as $\pi_{1}\left(B^{k b}\right)$. To complete the proof, we conside an element $A^{a} B^{5 b} . \pi_{k}$ is diagonal with eigenvalues $\beta^{5 k b}\left\{\alpha^{a}, \alpha^{3 a}, \alpha^{9 a}, \alpha^{5 a}, \alpha^{4 a}\right\}$. If $(a, 11)=1$ and $(b, 5)=1$ these are all primitive $55^{\text {th }}$ roots of unity; if $(a, 11)=1$ and $(b, 5)=5$ these are all primitive $11^{\text {th }}$ roots of unity; if $(a, 11)=11$ and $(b, 5)=1$ these are all primitive $5^{\text {th }}$ roots of unity. This shows $\pi_{k}\left(A^{a} B^{5 b}\right)$ is fixed point free and has the same eigenvalues as $\pi_{1}\left(A^{a} B^{5 k b}\right)$ which completes the proof.

We form the manifolds $M_{k}=S^{9} / \pi_{k}(G)$ with fundamental group $G$. These are all spherical space forms which inherit a natural orientation and metric from $S^{9}$ as discussed previously.

Lemma 4.9.6. Adopt the notation of Lemma 4.9.5. Let $M=S^{9} / \pi_{k}(G)$ be spherical space forms. Then $M_{1}, M_{2}, M_{3}$ and $M_{4}$ are all strongly isospectral.

Proof: Let $P$ be a self-adjoint elliptic differential operator which is natural in the category of oriented Riemannian manifolds. Let $P_{0}$ denote this operator on $S^{9}$, and let $P_{k}$ denote the corresponding operator on $M_{k}$. (For example, we could take $P$ to be the Laplacian on $p$-forms or to be the tangential operator of the signature complex). Let $\lambda \in \mathbf{R}$ and let $E_{0}(\lambda)$ and $E_{k}(\lambda)$ denote the eigenspaces of $P_{0}$ and $P_{k}$. We must show $\operatorname{dim} E_{k}(\lambda)$ is independent of $k$ for $1 \leq k \leq 4$. The unitary group acts on $S^{9}$ by orientation preserving isometries. The assumption of naturality lets us extend this to an action we shall denote by $e_{0}(\lambda)$ on $E_{0}(\lambda)$. Again, the assumption of naturality implies the eigenspace $E_{k}(\lambda)$ is just the subspace of $E_{0}(\lambda)$ invariant under the action of $e_{0}(\lambda)\left(\pi_{k}(G)\right)$. We can calculate the dimension
of this invariant subspace by:

$$
\begin{aligned}
\operatorname{dim} e_{k}(\lambda) & =\frac{1}{|G|} \sum_{g \in G} \operatorname{Tr}\left\{e_{0}(\lambda)\left(\pi_{k}(g)\right\}\right. \\
& =\frac{1}{275} \sum_{a, b} \operatorname{Tr}\left\{e_{0}(\lambda)\left(\pi_{k}\left(A^{a} B^{b}\right)\right\}\right.
\end{aligned}
$$

We apply Lemma 4.9 .5 to conclude $\pi_{k}\left(A^{a} B^{b}\right)$ is conjugate to $\pi_{1}\left(A^{a} B^{k b}\right)$ in $\mathrm{U}(5)$ so the two traces are the same and

$$
\begin{aligned}
\operatorname{dim} e_{k}(\lambda) & =\frac{1}{275} \sum_{a, b} \operatorname{Tr}\left\{e_{0}(\lambda) \pi_{1}\left(A^{a} B^{k b}\right)\right\} \\
& =\frac{1}{275} \sum_{a, b} \operatorname{Tr}\left\{e_{0}(\lambda) \pi_{1}\left(A^{a} B^{b}\right)\right\}=\operatorname{dim} e_{1}(\lambda)
\end{aligned}
$$

since we are just reparameterizing the group. This completes the proof.
Let $\rho \in R_{0}\left(\mathbf{Z}_{5}\right)$. We regard $\rho \in R_{0}(G)$ by defining $\rho\left(A^{a} B^{b}\right)=\rho(b)$. This is nothing but the pull-back of $\rho$ using the natural map $0 \rightarrow \mathbf{Z}_{55} \rightarrow G \rightarrow$ $\mathbf{Z}_{5} \rightarrow 0$.
Lemma 4.9.7. Let $\rho \in R_{0}\left(\mathbf{Z}_{5}\right)$ and let $G$ be as in Lemma 4.9.5. Let $\psi: G \rightarrow G$ be a group automorphism. Then $\rho \cdot \psi=\rho$ so this representation of $G$ is independent of the marking chosen.
Proof: The Sylow 11-subgroup is normal and hence unique. Thus $\psi(A)$ $=A^{a}$ for some $(a, 11)=1$ as $A$ generates the Sylow 11-subgroup. Let $\psi(B)=A^{c} B^{d}$ and compute:

$$
\psi\left(A^{3}\right)=A^{3 a}=\psi(B) \psi(A) \psi(B)^{-1}=A^{c} B^{d} A^{a} B^{-d} A^{-c}=A^{a \cdot 3^{d}}
$$

Since $(a, 11)=1, \quad 3^{d} \equiv 1$ (11). This implies $d \equiv 1$ (5). Therefore $\psi\left(A^{u} B^{v}\right)=A^{*} B^{d v}$ so that $\rho \psi\left(A^{u} B^{v}\right)=\rho(d v)=\rho(v)$ as $\rho \in R_{0}\left(\mathbf{Z}_{5}\right)$ which completes the proof.

These representations are canonical; they do not depend on the marking of the fundamental group. This defines a virtual locally flat bundle $V_{\rho}$ over each of the $M_{k}$. Let $P$ be the tangential operator of the signature complex; $\operatorname{ind}\left(\rho, \operatorname{signature}, M_{k}\right)$ is an oriented diffemorphism of $M_{k}$. In fact more is true. There is a canonical $\mathbf{Z}_{5}$ bundle over $M_{k}$ corresponding to the sequence $G \rightarrow \mathbf{Z}_{5} \rightarrow 0$ which by Lemma 4.9.7 is independent of the particular isomorphism of $\pi_{1}\left(M_{k}\right)$ with $G$ chosen. Lemma 4.6.4(b) shows this is a $\mathbf{Z}_{5}$-cobordism invariant. We apply Lemma 4.6.3 and calculate for $\rho \in R_{0}\left(\mathbf{Z}_{5}\right)$ that:
$\operatorname{ind}\left(\rho\right.$, signature, $\left.M_{k}\right)=\frac{1}{275} \cdot \sum_{a, b}^{\prime} \operatorname{Tr}(\rho(b)) \cdot \operatorname{defect}\left(\pi_{k}\left(A^{a} B^{b}\right)\right.$, signature $)$.
$\sum^{\prime}$ denotes the sum over $0 \leq a<11,0 \leq b<25,(a, b) \neq(0,0)$. Since this is an element of $R_{0}\left(\mathbf{Z}_{5}\right)$, we may suppose $(5, b)=1$. As $A^{a} B^{b}$ is conjugate to $B^{b}$ by 4.9.4, we can group the 11 equal terms together to see

$$
\operatorname{ind}\left(\rho, \text { signature, } M_{k}\right)=\frac{1}{25} \sum_{b}^{\prime} \operatorname{Tr}(\rho(b)) \cdot \operatorname{defect}\left(\pi_{k}\left(B^{b}\right), \text { signature }\right)
$$

where we sum over $0 \leq b<25$ and $(5, b)=1$.
If $(5, b)=1$, then the eigenvalues of $\pi_{k}\left(B^{b}\right)$ are $\left\{\beta^{b k} \varepsilon_{1}, \ldots, \beta^{b k} \varepsilon_{5}\right\}$. Thus

$$
\begin{aligned}
\operatorname{defect}\left(\pi_{k}\left(B^{b}\right), \text { signature }\right) & =\prod_{\nu} \frac{\beta^{b k} \varepsilon_{\nu}+1}{\beta^{b k} \varepsilon_{\nu}-1} \\
& =\frac{\beta^{5 b k}+1}{\beta^{5 b k}-1}
\end{aligned}
$$

since the product ranges over the primitive $5^{\text {th }}$ roots of unity. Let $\gamma=$ $\beta^{5}=e^{2 \pi i / 5}$, then we conclude

$$
\operatorname{ind}\left(\rho, \text { signature } \begin{array}{rl}
\left.M_{k}\right) & =\frac{1}{25} \cdot \sum_{b}^{\prime} \operatorname{Tr}(\rho(b)) \cdot \frac{\gamma^{k b}+1}{\gamma^{k b}-1} \\
& =\frac{1}{5} \cdot \sum_{0<b<5} \operatorname{Tr}(\rho(b)) \cdot \frac{\gamma^{k b}+1}{\gamma^{k b}-1}
\end{array}\right.
$$

if we group equal terms together.
We now calculate for the specific example $\rho=\rho_{1}-\rho_{0}$ :

$$
\operatorname{ind}\left(\rho_{1}-\rho_{0}, \text { signature, } \begin{array}{rl}
\left.M_{k}\right) & =\frac{1}{5} \sum_{b=1}^{4}\left(\gamma^{b}-1\right) \cdot \frac{\gamma^{k b}+1}{\gamma^{k b}-1} \\
& =\frac{1}{5} \sum_{b=1}^{4}\left(\gamma^{b \bar{k}}-1\right) \cdot \frac{\gamma^{b}+1}{\gamma^{b}-1}
\end{array}\right.
$$

if we let $k \bar{k} \equiv 1(5)$. We perform the indicated division; $\left(x^{\bar{k}}-1\right) /(x-1)=$ $x^{\bar{k}-1}+\cdots+1$ so we obtain

$$
=\frac{1}{5} \sum_{b=1}^{4}\left(\gamma^{b \bar{k}-b}+\gamma^{b \bar{k}-2 b}+\cdots+1\right)\left(\gamma^{b}+1\right)
$$

This expression is well defined even if $b=0$. If we sum over the entire group, we get an integer by the orthogonality relations. The value at 0 is $+2 \bar{k} / 5$ and therefore $\operatorname{ind}\left(\rho_{1}-\rho_{0}\right.$, signature, $\left.M_{k}\right)=-2 \bar{k} / 5 \bmod \mathbf{Z}$.

We choose the orientation arising from the given orientation on $S^{9}$. If we reverse the orientation, we change the sign of the tangential operator of the signature complex which changes the sign of this invariant.

Theorem 4.9.8. Let $M_{k}=S^{9} / \pi_{k}(G)$ with the notation of Lemma 4.9.5. These four manifolds are all strongly isospectral for $k=1,2,3,4$. There is a natural $\mathbf{Z}_{5}$ bundle over each $M_{k}$. ind $\left(\rho_{1}-\rho_{0}\right.$, signature, $\left.M_{k}\right)=-2 \bar{k} / 5$ $\bmod \mathbf{Z}$ where $k \bar{k} \equiv 1$ (5). As these 5 values are all different, these 4 manifolds are not $\mathbf{Z}_{5}$ oriented equivariant cobordant. Thus in particular there is no orientation preserving diffeomorphism between any two of these manifolds. $\quad M_{1}=-M_{4}$ and $M_{2}=-M_{3}$. There is no diffeomorphism between $M_{1}$ and $M_{2}$ so these are different topological types.
Proof: If we replace $\beta$ by $\bar{\beta}$ we replace $\pi_{k}$ by $\pi_{5-k}$ up to unitary equivalence. The map $z \mapsto \bar{z}$ reverses the orientation of $\mathbf{C}^{5}$ and thus $M_{k}=$ $-M_{5-k}$ as oriented Riemannian manifolds. This change just alters the sign of ind $(*, \operatorname{signature}, *)$. Thus the given calculation shows $M_{1} \neq \pm M_{2}$. The statement about oriented equivariant cobordism follows from section 4.6. (In particular, these manifolds are not oriented $G$-cobordant either.) This gives an example of strongly isospectral manifolds which are of different topological types.
Remark: These examples, of course, generalize; we have chosen to work with a particular example in dimension 9 to simplify the calculations involved.

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