# Sparse Approximations with Interior Point Methods 

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#### Abstract

Large-scale optimization problems that seek sparse solutions have become ubiquitous. They are routinely solved with various specialized first-order methods. Although such methods are often fast, they usually struggle with not-so-well conditioned problems. In this paper, specialized variants of an interior point-proximal method of multipliers are proposed and analyzed for problems of this class. Computational experience on a variety of problems, namely, multiperiod portfolio optimization, classification of data coming from functional Magnetic Resonance Imaging, restoration of images corrupted by Poisson noise, and classification via regularized logistic regression, provides substantial evidence that interior point methods, equipped with suitable linear algebra, can offer a noticeable advantage over first-order approaches.


## 1 Introduction

We are concerned with the efficient solution of a class of problems which are very large and are expected to yield sparse solutions. In practice, the sparsity is often induced by the presence of $\ell_{1}$ norm terms in the objective. We assume that a general problem of the following form

$$
\begin{array}{cl}
\min _{x}^{x} & f(x)+\tau_{1}\|x\|_{1}+\tau_{2}\|L x\|_{1}  \tag{1.1}\\
\text { s.t. } & A x=b,
\end{array}
$$

needs to be solved, where $f: \mathbb{R}^{n} \mapsto \mathbb{R}$ is a twice continuously differentiable convex function, $L \in \mathbb{R}^{l \times n}, A \in \mathbb{R}^{m \times n}, b \in \mathbb{R}^{m}$, and $\tau_{1}, \tau_{2}>0$. We are particularly interested in problems for which $f(x)$ displays some level of separability. The terms $\|x\|_{1}$ and $\|L x\|_{1}$ induce sparsity in the vector $x$ and/or in some (possibly redundant) dictionary $L x$. Numerous real-life problems can be recast into the form (1.1). Among
the various application areas, one can find portfolio optimization [40], signal and image processing [16, 50], classification in statistics [54] and machine learning [57], inverse problems [55] and compressed sensing [14], to mention just a few.

Optimization problems arising in these applications are usually solved by different specialized variants of first-order methods. Indeed, highly specialized and tuned to a narrow class of problems, first-order methods often outperform standard off-theshelf second-order techniques; the latter might be too expensive or might struggle with excessive memory requirements. Such comparisons are not fair though. With this paper we hope to change the incorrect opinion on the second-order methods. When efficiently implemented and specialized to a particular application, interior point methods (IPMs) offer an attractive alternative. They can be equally (or more) efficient than the best first-order methods available, and they deliver unmatched robustness and reliability.

The specializations of interior point methods proposed in this paper do not go beyond what has been commonly exploited by the first-order methods. Namely we propose:

- to exploit special features of the problems in the linear algebra of IPM, and
- to take advantage of the expected sparsity of the optimal solution.

In order to achieve our goals we propose to convert sparse approximation problems to standard smooth nonlinear convex programming ones by replacing the $\ell_{1}$ norm terms with a usual modeling trick in which an absolute value $|a|$ is substituted with the sum of two non-negative parts, $|a|=a^{+}+a^{-}$, where $a^{+}=\max \{a, 0\}$ and $a^{-}=\max \{-a, 0\}$. By introducing the auxiliary variable $d=L x \in \mathbb{R}^{l}$, problem (1.1) is then transformed into the following one:

$$
\begin{array}{cl}
\min _{x^{+}, x^{-}, d^{+}, d^{-}} & f\left(x^{+}-x^{-}\right)+\tau_{1}\left(e_{n}^{\top} x^{+}+e_{n}^{\top} x^{-}\right)+\tau_{2}\left(e_{l}^{\top} d^{+}+e_{l}^{\top} d^{-}\right), \\
\text {s.t. } & A\left(x^{+}-x^{-}\right)=b,  \tag{1.2}\\
& L\left(x^{+}-x^{-}\right)=d^{+}-d^{-}, \\
& x^{+}, x^{-}, d^{+}, d^{-} \geq 0,
\end{array}
$$

where $x^{+}, x^{-} \in \mathbb{R}^{n}$ are such that $x=x^{+}-x^{-}, d^{+}, d^{-} \in \mathbb{R}^{l}$ are such that $d=d^{+}-d^{-}$, and $e_{n} \in \mathbb{R}^{n}, e_{l} \in \mathbb{R}^{l}$ are vectors with entries equal to 1 . It is worth observing that (1.1) and (1.2) are equivalent; indeed the presence of linear terms which penalize for the sum of positive and negative parts of vectors $x$ and $d$ guarantee that at optimality only one of the split variables can take a nonzero value. Although (1.2) is larger than (1.1) because the variables have been replicated and new constraints have been added, it is in a form eligible to a direct application of an Interior Point Method (IPM). We expect that the well-known ability of IPMs to handle large sets of linear equality and non-negativity constraints will compensate for this increase of the problem dimension.

IPMs employ Newton method to solve a sequence of logarithmic barrier subproblems. In standard implementations of IPMs this requires many involved linear algebra operations (building and inverting the Hessian matrix), and for large problems it might become prohibitively expensive. In this paper we demonstrate that the use of inexact Newton method $[7,13,33]$ combined with a knowledgeable choice and appropriate tuning of linear algebra solvers (see [23, 26, 32, 35] and the references
therein) is the key to success when developing an IPM specialized to a particular class of problems. We also demonstrate an attractive ability of IPMs to select important variables and prematurely drop the irrelevant ones, a feature which is very well suited to solving sparse approximation problems in which the majority of variables are expected to be zero at optimality. It is worth mentioning at this point that our understanding of the features of IPMs applied to sparse approximation problems benefitted from the earlier studies which focused on compressed sensing problems [29, 30].

Ultimately, we provide computational evidence that IPMs can be more efficient than methods which are routinely used for the solution of sparse approximation problems by exploiting only first-order information.

## Notation

Throughout this paper we use lowercase Roman and Greek fonts to indicate scalars and vectors (the nature is clear from the context). Capital italicized Roman fonts are used to indicate matrices. Superscripts are used to denote the components of a vector/matrix. As an example, given $M \in \mathbb{R}^{m \times n}, v \in \mathbb{R}^{n} \mathcal{R} \subseteq\{1, \ldots, m\}$, and $\mathcal{C} \subseteq\{1, \ldots, n\}$, we set $v^{\mathcal{C}}:=\left(v^{i}\right)_{i \in \mathcal{C}}$ and $M^{\mathcal{R}, \mathcal{C}}:=\left(m^{i j}\right)_{i \in \mathcal{R}, j \in \mathcal{C}}$, where $v^{i}$ is the $i$-th entry of $v$ and $m^{i j}$ the $(i, j)$-th entry of $M$. We use $\lambda_{\min }(B)\left(\lambda_{\max }(B)\right.$, respectively) to denote the minimum (maximum) eigenvalue of an arbitrary square matrix $B$ with real eigenvalues. Similarly, $\sigma_{\min }(B)\left(\sigma_{\max }(B)\right.$, respectively) denotes the minimum (maximum) singular value of an arbitrary rectangular matrix $B$. We use $e_{n}$ and $0_{n}$ to denote a column vector of size $n$ with entries equal to 1 or 0 , respectively. Moreover, we use $I_{n}$ to indicate the identity matrix of size $n$ and $0_{m, n}$ to denote the zero matrix of size $m \times n$. We use subscripts to denote the elements of a sequence, e.g., $\left\{x_{k}\right\}$. Norms $\|\cdot\|$ are $\ell_{2}$, unless stated otherwise. For any finite set $\mathcal{A}$, we denote by $|\mathcal{A}|$ its cardinality. Finally, when referring to convex programming problems, we implicitly assume that the problems are linearly constrained.

## Structure of the article

The rest of this article is organized as follows. In Section 2 we briefly describe IPMs for convex programming, focusing in particular on the Interior Point-Proximal Method of Multipliers (IP-PMM), which is used in the subsequent sections. The choice of IP-PMM is motivated by the fact that it merges an infeasible IPM with the PMM approach, in order to keep the fast and reliable convergence properties of IPMs and the strong convexity of the PMM sub-problems, thus achieving better efficiency and robustness than both methods. In this section we also outline the testing environment used throughout the paper. In Sections 3 to 6 we present four applications formulated as optimization problems with sparsity sought in the solutions, and recast them in the form (1.2). In detail, in Section 3 we focus on a multi-period portfolio selection strategy, in Section 4 on the classification of data coming from functional Magnetic Resonance Imaging, in Section 5 on the restoration of images corrupted by Poisson noise, and in Section 6 on linear classification through regularized logistic regression. The first two applications yield convex quadratic programming problems, while the remaining ones yield general nonlinear convex programming problems. For each application, we provide a brief description of its
mathematical model and explain how IP-PMM is specialized for that case in terms of linear algebra solvers, including variable dropping strategies to help sparsification; we also show the results of computational experiments, including comparisons with state-of-the art methods widely used by the scientific community on the selected problems.

## 2 Interior Point Methods for Convex Programming

We consider the following convex programming problem:

$$
\begin{equation*}
\min _{x} f(x), \quad \text { s.t. } A x=b, \quad x \geq 0, \tag{2.1}
\end{equation*}
$$

where $x \in \mathbb{R}^{n}, A \in \mathbb{R}^{m \times n}$, and $f: \mathbb{R}^{n} \mapsto \mathbb{R}$ is a twice differentiable convex function. Using the Lagrangian duality theory [10], and defining a function $F(w): \mathbb{R}^{2 n+m} \mapsto$ $\mathbb{R}^{2 n+m}$, we write the KKT (optimality) conditions as follows:

$$
F(w)=\left[\begin{array}{c}
\nabla f(x)-A^{\top} y-z  \tag{2.2}\\
A x-b \\
X Z e_{n}
\end{array}\right]=\left[\begin{array}{c}
0_{n} \\
0_{m} \\
0_{n}
\end{array}\right],
$$

where $y \in \mathbb{R}^{m}$ and $z \in \mathbb{R}^{n}$ are the Lagrange multipliers corresponding to the equality and inequality constraints respectively, while $X, Z \in \mathbb{R}^{n \times n}$ denote the diagonal matrices with diagonal entries $x^{i}$ and $z^{i}$ (resp.), $\forall i \in\{1, \ldots, n\}$.

Problem (2.1) can be solved using a primal-dual IPM. There are numerous variants of IPMs and the reader is referred to [32] for an extended literature review. IPMs handle the non-negativity constraints of the problems with logarithmic barriers in the objective. That is, at each iteration $k$, we choose a barrier parameter $\mu_{k}$ and form the logarithmic barrier problem:

$$
\begin{equation*}
\min _{x} f(x)-\mu_{k} \sum_{j=1}^{n} \ln x^{j}, \quad \text { s.t. } A x=b . \tag{2.3}
\end{equation*}
$$

Then, Newton method (possibly an inexact variant of it [7, 33, 47]) is usually employed in order to approximately solve problem (2.3). Newton method is favored due to the self-concordance property of the logarithmic barrier [45]. Applying it to the optimality conditions of (2.3), and further forming the augmented system (by means of pivoting), we obtain a system of the following form:

$$
\left[\begin{array}{cc}
-\left(\nabla^{2} f\left(x_{k}\right)+\Theta_{k}^{-1}\right) & A^{\top}  \tag{2.4}\\
A & 0_{m, n}
\end{array}\right]\left[\begin{array}{c}
\Delta x_{k} \\
\Delta y_{k}
\end{array}\right]=\left[\begin{array}{c}
\nabla f\left(x_{k}\right)-A^{\top} y_{k}-\sigma_{k} \mu_{k} X_{k}^{-1} e \\
b-A x_{k}
\end{array}\right],
$$

where $\Theta_{k}=X_{k} Z_{k}^{-1}$. One can observe that the matrix $\Theta_{k}$ contains some very large and some very small elements close to optimality. Hence, the matrix in (2.4) becomes increasingly ill-conditioned, as the method progresses. Notice that as $\mu_{k} \rightarrow 0$, the optimal solution of (2.3) converges to the optimal solution of (2.1). Polynomial convergence of such methods, for various classes of problems, has been proved multiple times in the literature (see for example [45, 63]).

System like (2.4) can either be solved "exactly" (using an appropriate factorization, as in $[1,31,48]$ ) or iteratively (using an appropriate Krylov subspace method, as
in $[8,12,23,35]$ ). While the former approach is very general, it becomes problematic as the problem size increases. On the other hand, iterative methods (accompanied by appropriate preconditioners) may be difficult to generalize. However, if applied to specific classes of problems, they can allow one to solve huge-scale instances, by avoiding the explicit storage of the problem matrices.

### 2.1 Regularization in IPMs

In the context of IPMs, it is often beneficial to include some regularization, in order to improve the spectral properties of the system matrix in (2.4). For example, notice that if the constraint matrix $A$ is rank-deficient, then the matrix in (2.4) might not be invertible. The latter can be immediately addressed by the introduction of a dual regularization, say $\delta>0$, ensuring that $\operatorname{rank}\left(\left[A \delta I_{m}\right]\right)=m$. The introduction of a primal regularization, say $\rho>0$, can ensure that the matrix $\nabla^{2} f\left(x_{k}\right)+\Theta_{k}^{-1}+\rho I_{n}$, has eigenvalues that are bounded away from zero, and hence a significantly better worst-case conditioning than that of $\nabla^{2} f\left(x_{k}\right)+\Theta_{k}^{-1}$. To produce a diagonal term in the $(2,2)$ block Vanderbei added artificial variables to all the constraints [56]. Saunders and Tomlin $[51,52]$ achieved a similar result for the $(1,1)$ and $(2,2)$ blocks, by adding Tikhonov-type regularization terms to the original problem.

In later works, these Tikhonov-type regularization methods were replaced by algorithmic regularization schemes. In particular, one can observe that a very natural way of introducing primal regularization to problem (2.1), is through the application of the primal proximal point method. Similarly, dual regularization can be incorporated through the application of the dual proximal point method. This is a well-known fact. The authors in [1] presented a primal-dual regularized IPM for convex Quadratic Programming (QP), and interpreted this regularization as the application of the proximal point method. Subsequently, the authors in [31] developed a primal-dual regularized IPM, which applies PMM to solve convex QP problems, and employs a single IPM iteration for approximating the solution of each PMM sub-problem. There, global convergence of the method was proved, under some assumptions. A variation of the method proposed in [31] is given in [46], where general non-diagonal regularization matrices are employed, as a means of further improving factorization efficiency. Finally, the authors in [48] proposed an IP-PMM and proved (under standard assumptions) that it achieves convergence to an $\epsilon$-optimal solution in a polynomial number of iterations for convex QP problems and for linear semidefinite programming problems (see [47]), including cases when the associated Newton systems are solved inexactly.

### 2.2 Interior Point-Proximal Method of Multipliers

In this subsection, we derive an IP-PMM suitable for solving convex programming problems. The method is based on the developments in [48]. We consider the following primal problem (which can be equivalently formulated as (2.1), by adding some additional constraints):

$$
\begin{equation*}
\min _{x} f(x), \quad \text { s.t. } \quad A x=b, \quad x_{\mathcal{I}} \geq 0, \quad x_{\mathcal{F}} \text { free, } \tag{2.5}
\end{equation*}
$$

where $\mathcal{I} \subseteq\{1, \ldots, n\}, \mathcal{F}=\{1, \ldots, n\} \backslash \mathcal{I}$. In the above problem, we assume that the dimensions of the involved matrices are the same as those in (2.1). Effectively,
an IP-PMM arises by merging PMM with an infeasible IPM. For that purpose, assume that, at some iteration $k$ of the method, we have available an estimate $\eta_{k}$ for an optimal Lagrange multiplier vector $y^{*}$ associated to the equality constraints of (2.5). Similarly, we denote by $\zeta_{k}$ an estimate of a primal solution $x^{*}$. Now, we define the proximal penalty function that has to be minimized at the $k$-th iteration of the PMM, for solving (2.5), given the estimates $\eta_{k}, \zeta_{k}$ :

$$
\mathcal{L}_{\rho_{k}, \delta_{k}}^{P M M}\left(x ; \zeta_{k}, \eta_{k}\right)=f(x)-\eta_{k}^{\top}(A x-b)+\frac{1}{2 \delta_{k}}\|A x-b\|_{2}^{2}+\frac{\rho_{k}}{2}\left\|x-\zeta_{k}\right\|_{2}^{2}
$$

with $\left\{\delta_{k}\right\},\left\{\rho_{k}\right\}$ two positive non-increasing sequences of penalty parameters. Following [48], we require that these parameters decrease at the same rate as $\mu_{k}$; however, in practice we never allow these values to be reduced below a certain appropriately chosen threshold. For more details on how to choose these constants for general problems, the reader is referred to [46], where a perturbation analysis of regularization is performed. In order to solve the PMM sub-problem, we apply one (or a few) iterations of an infeasible IPM. To do that, we alter the previous penalty function, by including logarithmic barriers, that is

$$
\begin{equation*}
\mathcal{L}_{\rho_{k}, \delta_{k}}^{I P-P M M}\left(x ; \zeta_{k}, \eta_{k}\right)=\mathcal{L}_{\delta_{k}, \rho_{k}}^{P M M}\left(x ; \zeta_{k}, \eta_{k}\right)-\mu_{k} \sum_{j \in \mathcal{I}} \ln x^{j}, \tag{2.6}
\end{equation*}
$$

where $\mu_{k}>0$ is the barrier parameter. In order to form the optimality conditions of this sub-problem, we equate the gradient of $\mathcal{L}_{\delta_{k}, \rho_{k}}^{I P-P M M}$ with respect to $x$ to the zero vector, i.e.

$$
\nabla f(x)-A^{\top} \eta_{k}+\frac{1}{\delta_{k}} A^{\top}(A x-b)+\rho_{k}\left(x-\zeta_{k}\right)-\left[\begin{array}{c}
0_{|\mathcal{F}|} \\
\mu_{k}\left(X^{\mathcal{I}}\right)^{-1} e_{|\mathcal{I}|}
\end{array}\right]=0_{n}
$$

Next, we define the variables: $y=\eta_{k}-\frac{1}{\delta_{k}}(A x-b)$ and $z \in \mathbb{R}^{n}$, such that $z^{\mathcal{I}}=$ $\mu_{k}\left(X^{\mathcal{I}}\right)^{-1} e_{|\mathcal{I}|}, z^{\mathcal{F}}=0$, to get the following (equivalent) system of equations:

$$
\left[\begin{array}{c}
\nabla f(x)-A^{\top} y-z+\rho_{k}\left(x-\zeta_{k}\right) \\
A x+\delta_{k}\left(y-\eta_{k}\right)-b \\
X^{\mathcal{I}} z^{\mathcal{I}}-\mu_{k} e_{|\mathcal{I}|}
\end{array}\right]=\left[\begin{array}{c}
0_{n} \\
0_{m} \\
0_{n}
\end{array}\right]
$$

To solve the previous mildly nonlinear system of equations, at every iteration $k$, we employ Newton method and alter its right-hand side, using a centering parameter $\sigma_{k} \in(0,1)$. In other words, at every iteration of IP-PMM we have available an iteration triple $\left(x_{k}, y_{k}, z_{k}\right)$ and we want to solve the following system of equations:

$$
\begin{align*}
& {\left[\begin{array}{ccc}
-\left(\nabla^{2} f\left(x_{k}\right)+\rho_{k} I_{n}\right) & A^{\top} & I_{n} \\
A & \delta_{k} I_{m} & 0_{m, n} \\
Z_{k} & 0_{n, m} & X_{k}
\end{array}\right]\left[\begin{array}{c}
\Delta x \\
\Delta y \\
{\left[\begin{array}{c}
0_{|\mathcal{F}|} \\
\Delta z^{\mathcal{I}}
\end{array}\right]}
\end{array}\right] } \\
&=\left[\begin{array}{c}
\nabla f\left(x_{k}\right)-A^{\top} y_{k}+\sigma_{k} \rho_{k}\left(x_{k}-\zeta_{k}\right)-z_{k} \\
b-A x_{k}-\sigma_{k} \delta_{k}\left(y_{k}-\eta_{k}\right) \\
{\left[\begin{array}{c}
0_{|\mathcal{F}|} \\
\sigma_{k} \mu_{k} e_{|\mathcal{I}|}-X_{k}^{\mathcal{I}} z_{k}^{\mathcal{I}}
\end{array}\right]}
\end{array}\right] \tag{2.7}
\end{align*}
$$

From the third block-equation of (2.7) we have $\Delta z^{\mathcal{F}}=0$ and,

$$
\Delta z^{\mathcal{I}}=\left(X_{k}^{\mathcal{I}}\right)^{-1}\left(-Z_{k}^{\mathcal{I}} \Delta x^{\mathcal{I}}+\sigma_{k} \mu_{k} e_{|\mathcal{I}|}-X_{k}^{\mathcal{I}} z_{k}^{\mathcal{I}}\right)
$$

In light of the previous computations, (2.7) reduces to:

$$
\left[\begin{array}{cc}
-\left(\nabla^{2} f\left(x_{k}\right)+\Theta_{k}^{-1}+\rho_{k} I_{n}\right) & A^{\top}  \tag{2.8}\\
A & \delta_{k} I_{m}
\end{array}\right]\left[\begin{array}{l}
\Delta x \\
\Delta y
\end{array}\right]=\left[\begin{array}{l}
r_{1, k} \\
r_{2, k}
\end{array}\right]
$$

where

$$
\left[\begin{array}{c}
r_{1, k}  \tag{2.9}\\
r_{2, k}
\end{array}\right]=\left[\begin{array}{c}
\nabla f\left(x_{k}\right)-A^{\top} y_{k}+\sigma_{k} \rho_{k}\left(x_{k}-\zeta_{k}\right)-\left[\begin{array}{c}
0_{|\mathcal{F}|} \\
\sigma_{k} \mu_{k}\left(X_{k}^{\mathcal{I}}\right)^{-1} e_{|\mathcal{I}|}
\end{array}\right] \\
b-A x_{k}-\sigma_{k} \delta_{k}\left(y_{k}-\eta_{k}\right)
\end{array}\right]
$$

and

$$
\Theta_{k}^{-1}=\left[\begin{array}{cc}
0_{|\mathcal{F}|,|\mathcal{F}|} & 0_{|\mathcal{I}|,|\mathcal{F}|} \\
0_{|\mathcal{F}|,|\mathcal{I}|} & \left(X_{k}^{\mathcal{I}}\right)^{-1}\left(Z_{k}^{\mathcal{I}}\right)
\end{array}\right]
$$

In the special case where $\nabla^{2} f\left(x_{k}\right)$ is a diagonal (or zero) matrix, it could be beneficial to further reduce system (2.8), by eliminating variables $\Delta x$. The resulting normal equations yield a positive definite system of equations that reads as follows:

$$
\begin{align*}
& \left(A\left(\nabla^{2} f\left(x_{k}\right)+\Theta_{k}^{-1}+\rho_{k} I_{n}\right)^{-1} A^{\top}+\delta_{k} I_{m}\right) \Delta y  \tag{2.10}\\
= & r_{2, k}+A\left(\nabla^{2} f\left(x_{k}\right)+\Theta_{k}^{-1}+\rho_{k} I_{n}\right)^{-1}\left(r_{1, k}\right)
\end{align*}
$$

The parameters $\eta_{k}, \zeta_{k}$ are tuned as in [48]. In particular, we set $\eta_{0}=y_{0}$ and $\zeta_{0}=x_{0}$, where $\left(x_{0}, y_{0}, z_{0}\right)$ is the starting point of IP-PMM. Then, at the end of every iteration $k$, we set $\left(\zeta_{k+1}, \eta_{k+1}\right)=\left(x_{k+1}, y_{k+1}\right)$, only if the primal and dual residuals are decreased sufficiently. If the latter is not the case, we set $\left(\zeta_{k+1}, \eta_{k+1}\right)=\left(\zeta_{k}, \eta_{k}\right)$.

It has been demonstrated in [48] that IP-PMM using a single Newton step per iteration converges to an $\epsilon$-optimal solution in a polynomial number of iterations, if $f$ is a convex quadratic function. Furthermore, the latter holds for linear semi-definite programming problems, even if one solves the Newton system inexactly (see [47]). Nevertheless, the previous is not proven to hold for the general convex (nonlinear) case. In the latter case, one would have to employ Newton method combined with a line-search or a trust-region strategy (e.g. see [58]), in order to guarantee the convergence of the method. In all the cases analyzed in this work we make use of a simple Mehrotra-type (see [41]) predictor-corrector scheme approximating the central path, which in general is sufficient to produce good directions that allow the method to converge quickly to the optimal solution. In the corrector stage, the right-hand side is approximated by the linearization of the function that is being minimized (see [53]).

### 2.3 Testing environment

The various specialiations of IP-PMM discussed in the following sections have all been implemented in MATLAB and compared with MATLAB implementations of state-of-the-art methods for each specific problem. All the tests were run with MATLAB R2019b on an Intel Xeon Platinum 8168 CPU with 192 GB of RAM, available from the magicbox server at the Department of Mathematics and Physics of the University of Campania "L. Vanvitelli".

## 3 Portfolio Selection Problem

Portfolio selection is one of the most central topics in modern financial economics. It deals with the decision problem of how to allocate resources among several competing assets in accordance with the investment objectives. For medium- and long-time horizons, the multi-period strategy is suitable, because it allows the change of the allocation of the capital across the assets, taking into account the evolution of the available information. In a multi-period setting, the investment period is partitioned into $m$ sub-periods, delimited by $m+1$ rebalancing dates $t_{j}$. The decisions are taken at the beginning of each sub-period $\left[t_{j}, t_{j+1}\right), j=1, \ldots, m$, and kept within it. The optimal portfolio is defined by the vector

$$
w=\left[w_{1}^{\top}, w_{2}^{\top}, \ldots, w_{m}^{\top}\right]^{\top}
$$

where $w_{j} \in \mathbb{R}^{s}$ is the portfolio of holdings at time $t_{j}$ and $s$ is the number of assets.
The mean-variance formulation proposed by Markowitz [40] was extended to a multi-period portfolio selection by Li and Ng [38], and in recent years there has been a significant advancement of both theory and methodologies. In a multi-period mean-variance framework, we fix a final target expected return and adopt as risk measure the function obtained by summing the single-period variance terms [17]:

$$
\rho(w)=\sum_{j=1}^{m} w_{j}^{\top} C_{j} w_{j}
$$

where $C_{j} \in \mathbb{R}^{s \times s}$ is the covariance matrix, assumed to be positive definite, estimated at $t_{j}$. A common strategy to estimate Markowitz model parameters is to use historical data as predictive of the future behavior of asset returns. Different regularization techniques have been proposed to deal with ill-conditioning due to asset correlation; in the last years the $\ell_{1}$-regularization has been used to promote sparsity in the solution [21]. It allows investors to reduce the number of positions to be monitored and the holding and transaction costs. Another useful interpretation of the $\ell_{1}$ norm is related to the amount of short positions, which indicate an investment strategy where an investor is selling borrowed stocks in the open market, expecting that the market will drop, in order to realize a profit. A suitable tuning of the regularization parameter permits a short controlling in both the single- and the multi-period case [18, 21]. However, in the multi-period case, the sparsity in the solution does not guarantee the control of the transaction costs, especially if the pattern of the active positions completely changes across periods. In this case, sparsity must be introduced in the variation, e.g. by adding an $\ell_{1}$ term involving the differences of the wealth values allocated on the assets between two contiguous rebalancing times. This acts as a penalty on the portfolio turnover, which has the effect of reducing the number of transactions and hence the transaction costs [19, 24].

Thus, we consider the following fused lasso optimization problem for multi-period portfolio selection [24]:

$$
\begin{align*}
\min _{w} & \frac{1}{2} w^{\top} C w+\tau_{1}\|w\|_{1}+\tau_{2}\|L w\|_{1}, \\
\text { s.t. } & w_{1}^{\top} e_{s}=\xi_{\text {init }},  \tag{3.1}\\
& w_{j}^{\top} e_{s}=\left(e_{s}+r_{j-1}\right)^{\top} w_{j-1}, \quad j=2, \ldots, m, \\
& \left(e_{s}+r_{m}\right)^{\top} w_{m}=\xi_{\text {term }},
\end{align*}
$$

where $n=m s, C=\operatorname{diag}\left(C_{1}, C_{2}, \ldots, C_{m}\right) \in \mathbb{R}^{n \times n}$ is a block diagonal symmetric positive definite matrix, $\tau_{1}, \tau_{2}>0, L \in \mathbb{R}^{(n-s) \times n}$ is the discrete difference operator representing the fused lasso regularizer, $r_{j} \in \mathbb{R}^{s}$ is the expected return vector at time $t_{j}, \xi_{\text {init }}$ is the initial wealth, and $\xi_{\text {term }}$ is the target expected wealth resulting from the overall investment. The first constraint is the initial budget constraint. The strategy is assumed to be self-financing, as constraints from 2 to $m$ establish; this means that the value of the portfolio changes only because the asset prices change. The $(m+1)$-st constraint defines the expected final wealth. To deal with the nonseparability of the objective function in (3.1), we introduce an auxiliary variable $d$, which is constrained to be equal to $L w$, and we equivalently formulate problem (3.1) as follows:

$$
\begin{array}{ll}
\min _{w, d} & \frac{1}{2} w^{\top} C w+\tau_{1}\|w\|_{1}+\tau_{2}\|d\|_{1} \\
\text { s.t. } & \bar{A} w=\bar{b}  \tag{3.2}\\
& L w=d
\end{array}
$$

where the constraint matrix $\bar{A} \in \mathbb{R}^{(m+1) \times n}$ can be interpreted as a $(m+1) \times m$ lower bi-diagonal block matrix, with blocks of dimension $1 \times s$ defined by

$$
\bar{A}^{i, j}=\left\{\begin{array}{cl}
e_{s}^{\top} & \text { if } i=j \\
-\left(e_{s}+r_{i-1}\right)^{\top} & \text { if } j=i+1 \\
0_{s}^{\top} & \text { otherwise }
\end{array}\right.
$$

and $\bar{b}=\left(\xi_{\text {init }}, 0,0, \ldots, \xi_{\text {term }}\right)^{\top} \in \mathbb{R}^{m+1}$.

### 3.1 Specialized IP-PMM for quadratic portfolio optimization problems

Using the standard trick described in Section 1, we split $w$ and $d$ into two vectors of the same size, representing the non-negative and non-positive parts of the entries of $w$ and $d$ respectively, i.e., $w=w^{+}-w^{-}$and $d=d^{+}-d^{-}$. Then, problem (3.2) is reformulated as the following QP problem:

$$
\begin{equation*}
\min _{x} \frac{1}{2} x^{\top} Q x+c^{\top} x, \quad \text { s.t. } \quad A x=b, \quad x \geq 0 \tag{3.3}
\end{equation*}
$$

where we set $l=n-s, \bar{n}=2(n+l)=2 s(2 m-1), \bar{m}=m+1+l=(m+1)+s(m-1)$,

$$
\begin{gather*}
x=\left[\left(w^{+}\right)^{\top},\left(w^{-}\right)^{\top},\left(d^{+}\right)^{\top},\left(d^{-}\right)^{\top}\right]^{\top} \in \mathbb{R}^{\bar{n}}, \\
Q=\left[\begin{array}{rr}
{\left[\begin{array}{rr}
C & -C \\
-C & C
\end{array}\right]} & 0_{2 n, 2 l} \\
0_{2 l, 2 n} & 0_{2 l, 2 l}
\end{array}\right] \in \mathbb{R}^{\bar{n} \times \bar{n}}, \quad A=\left[\begin{array}{ccc}
\bar{A} & -\bar{A} & 0_{(m+1), 2 l} \\
L & -L & {\left[\begin{array}{ll}
-I_{l} & I_{l}
\end{array}\right]}
\end{array}\right] \in \mathbb{R}^{\bar{m} \times \bar{n}}, \\
c=\left[\tau_{1}, \ldots, \tau_{1}, \tau_{2}, \ldots, \tau_{2}\right]^{\top} \in \mathbb{R}^{\bar{n}}, \quad b=\left[\bar{b}^{1}, \ldots, \bar{b}^{m+1}, 0, \ldots, 0\right]^{\top} \in \mathbb{R}^{\bar{m}} . \tag{3.4}
\end{gather*}
$$

### 3.1.1 Dropping Primal Variables

The optimal solution of problem (3.2) is expected to be sparse. On the other hand, in light of the reformulation given in (3.3), we anticipate (and verify in practice) that most of the primal variables $w$ attain a zero value close to optimality. Such
variables may significantly contribute to the ill conditioning of the matrix in (2.8). In order to take advantage of this special property displayed by problem (3.3), we employ the following heuristic method, which aims at dropping variables $x^{j}$ which are sufficiently close zero, their seemingly optimal value. This results in better conditioning of the augmented system, whose dimension is also significantly reduced close to optimality, thus decreasing the computational cost of each IPM iteration. In other words, as IP-PMM progresses, we project the problem onto a smaller space, with the assumption that the projected (smaller) problem has the same optimal solution as that of the original problem.

In particular, we set a threshold value $\epsilon_{\text {drop }}>0$, and a large constant $\xi>0$. At iteration $k=0$, we define a set $\mathcal{V}=\emptyset$. Then, at every iteration $k$ of IP-PMM, we check the following condition, for every $j \in \mathcal{I} \backslash \mathcal{V}$ :

$$
\begin{equation*}
x_{k}^{j} \leq \epsilon_{\text {drop }} \quad \text { and } \quad z_{k}^{j} \geq \xi \cdot \epsilon_{\text {drop }} \quad \text { and } \quad\left(r_{d}\right)_{k}^{j} \leq \epsilon_{\text {drop }} \tag{3.5}
\end{equation*}
$$

where $\left(r_{d}\right)_{k}^{j}=\left(c-A^{\top} y_{k}+Q x_{k}-z_{k}\right)^{j}$ represents the dual infeasibility corresponding to the $j$-th variable. Any variable that satisfies the latter condition is dropped, that is, we set $x_{k}^{j}=0, \mathcal{V}=\mathcal{V} \cup\{j\}, \mathcal{G}=\mathcal{F} \cup(\mathcal{I} \backslash \mathcal{V})$, we drop $z_{k}^{j}$ and solve

$$
\left[\begin{array}{cc}
-\left(Q^{\mathcal{G}, \mathcal{G}}+\left(\Theta_{k}^{\mathcal{G}, \mathcal{G}}\right)^{-1}+\rho_{k} I_{|\mathcal{G}|}\right) & \left(A^{\mathcal{H}, \mathcal{G}}\right)^{\top}  \tag{3.6}\\
A^{\mathcal{H}, \mathcal{G}} & \delta_{k} I_{\bar{m}}
\end{array}\right]\left[\begin{array}{l}
\Delta x \\
\Delta y
\end{array}\right]=\left[\begin{array}{c}
r_{1, k}^{\mathcal{G}} \\
r_{2, k}
\end{array}\right],
$$

where, in the presence of free variables, $\left(\Theta_{k}^{\mathcal{F}, \mathcal{F}}\right)^{-1}$ is considered to be the zero matrix, $r_{1, k}, r_{2, k}$ are defined in (2.9) (by substituting $A^{\mathcal{H}, \mathcal{G}}$ as the constraint matrix), and $\mathcal{H}=\{1, \ldots, \bar{m}\}$. We should note that this is a heuristic, since once a variable is dropped, it is not considered again until the method converges. Hence, one has to make sure that none of the nonzero variables $x_{k}^{j}$ is dropped. Nevertheless, at the end of the optimization process we can test whether the variables in $\mathcal{V}$ are indeed nonzero. More specifically, once an optimal solution $\left(x^{*}, y^{*}, z^{*}\right)$ is found, we compute:

$$
z^{\mathcal{V}}=c^{\mathcal{V}}-\left(A^{\mathcal{H}, \mathcal{V}}\right)^{\top} y^{*}+Q^{\mathcal{V}, \mathcal{G}}\left(x^{*}\right)^{\mathcal{G}}
$$

If there exists $j$ such that $\left(z^{\mathcal{V}}\right)^{j} \leq 0$, then we would identify that a variable $x^{j}$ was incorrectly dropped. Otherwise, the optimal solution of the reduced problem coincides with the optimal solution of (3.3). We notice that this methodology is not new. In particular, a similar strategy was employed in [34], where a special class of linear programming problems was solved using a primal-dual barrier method.

### 3.2 Computational Experience

We test the effectiveness of the IP-PMM applied to the fused lasso model on the following real-world data sets:

1. FF48-FF100 (Fama \& French 48-100 Industry portfolios, USA), containing 48-100 portfolios considered as assets, from July 1926 to December 2015.
2. ES50 (EURO STOXX 50), containing 50 stocks from 9 Eurozone countries (Belgium, Finland, France, Germany, Ireland, Italy, Luxembourg, the Netherlands and Spain), from January 2008 to December 2013.
3. FTSE100 (Financial Times Stock Exchange, UK), containing 100 assets, from July 2002 to April 2016.
4. SP500 (Standard \& Poors, USA), containing 500 assets, from January 2008 to December 2016.
5. NASDAQC (National Association of Securities Dealers Automated Quotation Composite, USA), containing almost all stocks listed on the Nasdaq stock market, from February 2003 to April 2016.

Following [19, 20], we generate 10 problems with annual or quarterly rebalancing, after a preprocessing procedure that eliminates the elements with the highest volatilities. A rolling window (RW) for setting up the model parameters is considered. For each dataset, the length of the RWs is fixed in order to build positive definite covariance matrices and ensure statistical significance. Different datasets require different lengths for the RWs. In Table 1 we summarize the information on the test problems.

Table 1: Characteristics of the portfolio test problems ( $\mathrm{y}=$ years, $\mathrm{m}=$ months)

| Problem | Assets | RW | Sub-periods | $\bar{n}$ |
| :--- | ---: | ---: | ---: | ---: |
| FF48-10 | 48 | 5 y | 10 y | 1632 |
| FF48-20 | 48 | 5 y | 20 y | 3552 |
| FF48-30 | 48 | 5 y | 30 y | 5472 |
| FF100-10 | 96 | 10 y | 10 y | 3264 |
| FF100-20 | 96 | 10 y | 20 y | 7104 |
| FF100-30 | 96 | 10 y | 30 y | 10,944 |
| ES50 | 50 | 1 y | 22 m | 4300 |
| FTSE100 | 83 | 1 y | 10 y | 3154 |
| SP500 | 431 | 2 y | 8 y | 11,206 |
| NASDAQC | 1203 | 1 y | 10 y | 45,714 |

We introduce some measures to evaluate the goodness of the optimal portfolios versus the benchmark one, in terms of risk, sparsity and transaction costs. We consider as benchmark the multi-period naive portfolio, based on the strategy for which at each rebalancing date the total amount is equally divided among the assets. We assume that the investor has one unit of wealth at the beginning of the planning horizon, i.e., $\xi_{\text {init }}=1$, and we set as expected final wealth the one provided by the benchmark. As in [19], we define:

$$
\begin{equation*}
\text { ratio }=\frac{w_{\text {naive }}^{\top} C w_{\text {naive }}}{w_{\text {opt }}^{\top} C w_{\text {opt }}}, \tag{3.7}
\end{equation*}
$$

where $w_{\text {naive }}$ and $w_{o p t}$ are respectively the naive portfolio and the optimal one. This value measures the risk reduction factor with respect to the benchmark. We consider the number of active positions as a measure of holding costs; then the value

$$
\begin{equation*}
\text { ratio }_{h}=\frac{\text { \# active positions of } w_{\text {naive }}}{\# \text { active positions of } w_{o p t}} \tag{3.8}
\end{equation*}
$$

measures the reduction factor of the holding costs with respect to the benchmark. Finally, we consider the number of variations in the weights as a measure of transaction costs. More precisely, if $w_{j}^{i} \neq w_{j+1}^{i}$ we assume that security $i$ has been bought or sold in the period $\left[t_{j}, t_{j+1}\right)$, then we estimate the number of transactions as:

$$
\mathcal{T}=\operatorname{trace}\left(V^{\top} V\right)
$$

where $V \in \mathbb{R}^{s \times(m-1)}$, with

$$
v^{i j}= \begin{cases}1 & \text { if }\left|w_{j}^{i}-w_{j+1}^{i}\right| \geq \epsilon \\ 0 & \text { otherwise }\end{cases}
$$

and $\epsilon>0$, in order to make sense in financial terms. A measure of the transaction reduction factor with respect to the benchmark is given by

$$
\begin{equation*}
\text { ratio }_{t}=\frac{\mathcal{T}_{\text {naive }}}{\mathcal{T}_{\text {opt }}} \tag{3.9}
\end{equation*}
$$

We consider a version of the presented IP-PMM algorithm in which the solution of problem (3.6) is computed exactly, the parameter $\epsilon_{\text {drop }}$ controlling the heuristic described in Section 3.1.1 is set to $10^{-4}$, and the constant $\xi$, which is used to ensure that the respective dual slack variable is bounded away from zero, is set to $10^{2}$. We compare IP-PMM with the Split Bregman method, which is known to be very efficient for this kind of problems. In detail, we consider the Alternating Split Bregman algorithm used in [20], based on a further reformulation of problem (3.2) as

$$
\begin{aligned}
\min _{w, u, d} & \frac{1}{2} w^{\top} C w+\tau_{1}\|u\|_{1}+\tau_{2}\|d\|_{1}, \\
\text { s.t. } & \bar{A} w=\bar{b} \\
& L w=d \\
& w=u
\end{aligned}
$$

This algorithm splits the minimization in three parts. Given $w_{k}, u_{k}, d_{k}$, the $(k+1)$-st iteration consists in the minimization of a quadratic function to determine $w_{k+1}$ and the application of the soft-thresholding operator

$$
[\mathcal{S}(v, \gamma)]^{i}=\frac{v^{i}}{\left|v^{i}\right|} \max \left(\left|v^{i}\right|-\gamma, 0\right)
$$

where $v$ is a real vector and $\gamma>0$, to determine $u_{k+1}$ and $d_{k+1}$. The optimal value $w_{k+1}$ can be obtained by solving the system $H w=p_{k+1}$, where

$$
\begin{equation*}
H=C+\lambda_{1} A^{\top} A+\lambda_{2} L^{\top} L+\lambda_{3} I \tag{3.10}
\end{equation*}
$$

where $\lambda_{1}, \lambda_{2}, \lambda_{3}>0$ are fixed and $p_{k+1}$ depends on the iteration. Since $H$ is independent of the iteration and is symmetric positive definite and sparse, in [20] the authors compute its sparse Cholesky factorization only once and solve two triangular systems at each iteration. We refer to this algorithm as ASB-Chol.

The values of $\tau_{1}$ and $\tau_{2}$ in (3.1) are selected to guarantee reasonable portfolios in terms of short positions, i.e., in terms of negative components in the solution. Note that from the financial point of view, negative solutions correspond to transactions

Table 2: IP-PMM vs ASB-Chol

|  | IP-PMM |  |  |  |  |  |
| :--- | ---: | ---: | ---: | ---: | ---: | :---: |
| Problem | Time (s) | Iters | ratio | ratio | ratio |  |
| FF48-10 | $1.37 \mathrm{e}-1$ | 12 | $2.32 \mathrm{e}+0$ | $6.67 \mathrm{e}+0$ | $1.66 \mathrm{e}+1$ |  |
| FF48-20 | $3.77 \mathrm{e}-1$ | 16 | $2.28 \mathrm{e}+0$ | $6.58 \mathrm{e}+0$ | $2.13 \mathrm{e}+1$ |  |
| FF48-30 | $8.43 \mathrm{e}-1$ | 21 | $4.64 \mathrm{e}+0$ | $6.15 \mathrm{e}+0$ | $1.69 \mathrm{e}+1$ |  |
| FF100-10 | $4.92 \mathrm{e}-1$ | 12 | $1.58 \mathrm{e}+0$ | $1.78 \mathrm{e}+1$ | $4.36 \mathrm{e}+1$ |  |
| FF100-20 | $1.63 \mathrm{e}+0$ | 15 | $1.81 \mathrm{e}+0$ | $2.04 \mathrm{e}+1$ | $4.92 \mathrm{e}+1$ |  |
| FF100-30 | $3.93 \mathrm{e}+0$ | 21 | $5.82 \mathrm{e}+0$ | $1.34 \mathrm{e}+1$ | $3.60 \mathrm{e}+1$ |  |
| ES50 | $4.59 \mathrm{e}-1$ | 14 | $2.12 \mathrm{e}+0$ | $4.42 \mathrm{e}+0$ | $5.75 \mathrm{e}+1$ |  |
| FTSE100 | $4.64 \mathrm{e}-1$ | 14 | $1.85 \mathrm{e}+0$ | $5.37 \mathrm{e}+1$ | $6.09 \mathrm{e}+1$ |  |
| SP500 | $3.43 \mathrm{e}+1$ | 16 | $1.57 \mathrm{e}+0$ | $8.62 \mathrm{e}+1$ | $1.50 \mathrm{e}+2$ |  |
| NASDAQC | $7.05 \mathrm{e}+2$ | 20 | $3.15 \mathrm{e}+0$ | $2.73 \mathrm{e}+0$ | $3.89 \mathrm{e}+2$ |  |
|  |  |  |  |  |  |  |
| ASB-Chol |  |  |  |  |  |  |
| Problem | Time (s) | Iters | ratio | ratio | ratio |  |
| FF48-10 | $1.67 \mathrm{e}-1$ | 1431 | $2.33 \mathrm{e}+0$ | $6.67 \mathrm{e}+0$ | $1.66 \mathrm{e}+1$ |  |
| FF48-20 | $3.72 \mathrm{e}-1$ | 1985 | $2.31 \mathrm{e}+0$ | $7.93 \mathrm{e}+0$ | $2.09 \mathrm{e}+1$ |  |
| FF48-30 | $1.12 \mathrm{e}+0$ | 4125 | $4.64 \mathrm{e}+0$ | $6.08 \mathrm{e}+0$ | $1.66 \mathrm{e}+1$ |  |
| FF100-10 | $8.49 \mathrm{e}-1$ | 3087 | $1.58 \mathrm{e}+0$ | $1.78 \mathrm{e}+1$ | $4.36 \mathrm{e}+1$ |  |
| FF100-20 | $2.09 \mathrm{e}+0$ | 3635 | $1.80 \mathrm{e}+0$ | $1.78 \mathrm{e}+1$ | $4.27 \mathrm{e}+1$ |  |
| FF100-30 | $8.54 \mathrm{e}+0$ | 9043 | $5.83 \mathrm{e}+0$ | $1.12 \mathrm{e}+1$ | $2.97 \mathrm{e}+1$ |  |
| ES50 | $9.70 \mathrm{e}-1$ | 4297 | $2.05 \mathrm{e}+0$ | $2.94 \mathrm{e}+0$ | $4.26 \mathrm{e}+1$ |  |
| FTSE100 | $4.29 \mathrm{e}-1$ | 1749 | $1.80 \mathrm{e}+0$ | $5.07 \mathrm{e}+1$ | $5.71 \mathrm{e}+1$ |  |
| SP500 | $1.98 \mathrm{e}+1$ | 3728 | $1.74 \mathrm{e}+0$ | $6.16 \mathrm{e}+1$ | $1.01 \mathrm{e}+2$ |  |
| NASDAQC | $8.84 \mathrm{e}+2$ | 14264 | $3.15 \mathrm{e}+0$ | $2.73 \mathrm{e}+0$ | $3.89 \mathrm{e}+2$ |  |

in which an investor sells borrowed securities in anticipation of a price decline. In our runs we consider the smallest values of $\tau_{1}$ and $\tau_{2}$ that produce at most $2 \%$ of short positions. We set $\tau_{1}=\tau_{2}=10^{-2}$ for the FF48 and FF100 data sets, $\tau_{1}=\tau_{2}=10^{-3}$ for ES50 and SP500, $\tau_{1}=10^{-2}$ and $\tau_{2}=10^{-3}$ for FTSE, and $\tau_{1}=10^{-2}$ and $\tau_{2}=10^{-4}$ for NASDAQC.

In Table 2 we present the results obtained with IP-PMM and ASB-Chol on the test problems. The termination criteria of IP-PMM are the same as in [48], while those of ASB-Chol are the usual ones based on the violation of constraints. The relative tolerance for the two algorithms is $t o l=10^{-6}$, which guarantees that the values of ratio differ by at most $10 \%$, so that both algorithms produce comparable portfolios in terms of risk. We note that the solution computed by ASB-Chol is thresholded by setting to zero all the entries with absolute value not exceeding the same value of $\epsilon_{\text {drop }}$ used in the IP-PMM dropping strategy. The results show that the optimal portfolios computed by IP-PMM and ASB-Chol outperform the benchmark ones in terms of all the metrics. Concerning ratio and ratio , IP-PMM is generally able to produce greater values than ASB-Chol, which indicates a higher sparsity in the solution found by IP-PMM. IP-PMM generally performs comparably or better than ASB-Chol in terms of elapsed time. Although ASB-Chol is faster that IP-PMM on SP500 by 14.5 seconds ( $42 \%$ ), IP-PMM is able to reach a better solution in terms of sparsity. When applied to FF100-30 IP-PMM produces a portfolio associated with
lower transaction costs and takes less than half of the time required by ABS-Chol. When applied to NASDAQC, which is the largest problem under consideration, the two methods reach comparable solutions in terms of all the evaluation metrics, but IP-PMM needs about $20 \%$ less time ( 179 seconds) than ASB-Chol. This suggests that the use of IP-PMM can be beneficial especially when solving high-dimensional problems.

## 4 Classification models for functional Magnetic Resonance Imaging data

The functional Magnetic Resonance Imaging (fMRI) technique measures brain spatiotemporal activity via Blood-Oxygen-Level-Dependent (BOLD) signals. Starting from the assumption that neuronal activity is coupled with cerebral blood flow, fMRI signals have been used to identify regions associated with functions such as speaking, vision, movement, etc.. By analyzing the different oxigenation levels in specific areas of the brain of healthy and ill patients, in the last decades fMRI signals have been used to investigate the effect on the brain functionality of tumors, strokes, head and brain injuries and of cognitive disorders such as schizophrenia or Alzheimer's (see [27, 37, 42] and the references therein).

In an fMRI scan, voxels representing regions of the brain of a patient are recorded at different time instances. The temporal resolution is usually in the order of a few seconds, while the spatial resolution generally ranges from $4-5 \mathrm{~mm}$ (for some full brain analyses) to 1 mm (for analyses on specific brain regions), which may amount up to about a million voxels. Since fMRI experiments are conducted over groups of patients, the dimensionality of the data is further increased. Therefore, the interpretation of fMRI results requires the analysis of huge quantities of data. To this aim, machine learning techniques are being increasingly used in recent years, because of their capability of dealing with massive amounts of data, incorporating also apriori information about the problems they are targeted to [3, 4, 27, 28, 36, 42, 49].

Here we focus on the problem of training a binary linear classifier to distinguish between different classes of patients (e.g., ill/healthy) or different kinds of stimuli (e.g., pleasant/unpleasant), and to get information about the most significant brain areas associated with the related neural activity. The two classes are identified by the labels -1 and 1 . We assume that the training set consists of $s_{-1} 3$-dimensional (3d) scans in class -1 and $s_{1} 3 \mathrm{~d}$ scans in class 1 , where each 3 d scan is reshaped as a row vector of size $q=q_{1} \times q_{2} \times q_{3}$, and $q_{i}$ is the number of voxels along the $i$-th coordinate direction of the domain covering the brain. All the scans are stored as rows of a matrix $D \in \mathbb{R}^{s \times q}$, where $s=s_{-1}+s_{1}$.

Following [4], we use a square loss function with the aim of determining an unbiased hyperplane in $\mathbb{R}^{q}$ that can separate the patients in the two classes. This leads to a minimization problem of the form:

$$
\begin{equation*}
\min \frac{1}{2 s}\|D w-\hat{y}\|^{2} \tag{4.1}
\end{equation*}
$$

where $\hat{y}$ is a vector containing the labels associated with each scan. Since the number of patients is usually much smaller than the size of a scan, i.e., $s \ll q$, problem (4.1) is strongly ill posed and thus requires regularization. Recently, significant attention has
been given to regularization terms encouraging the presence of structured sparsity, where smoothly varying non-zero coefficients of the solution are associated with small contiguous regions of the brain. This is motivated by the possibility of obtaining more interpretable solutions than those corresponding to other regularizers that do not promote sparsity or lead to sparse solutions without any structure (see [4] and the references therein).

Structured sparsity can be promoted, e.g., by using a combination of $\ell_{1}$ and anisotropic Total Variation (TV) terms [2], which can be regarded as a fused lasso regularizer [54]. The regularized problem reads

$$
\begin{equation*}
\min _{w} \frac{1}{2 s}\|D w-\hat{y}\|^{2}+\tau_{1}\|w\|_{1}+\tau_{2}\|L w\|_{1} \tag{4.2}
\end{equation*}
$$

where $\|L w\|_{1}$ is the discrete anisotropic TV of $w$, i.e., $L=\left[\begin{array}{lll}L_{x}^{\top} & L_{y}^{\top} & L_{z}^{\top}\end{array}\right]^{\top} \in \mathbb{R}^{l \times q}$ is the matrix representing first-order forward finite differences in the $x, y, z$-directions at each voxel. By penalizing the difference between each voxel and its neighbors in each direction, one enforces the weights of the classification hyperplane (which share the 3d structure of the scans) to assume similar values for contiguous regions of the brain, thus leading to identify whole regions of the brain involved in the decision process.

The previous problem can be reformulated by introducing the variables $u=D w$ and $d=L w$ and applying the splitting

$$
w=w^{+}-w^{-}, d=d^{+}-d^{-}, \quad\left(w^{+}, w^{-}, d^{+}, d^{-}\right) \geq\left(0_{q}, 0_{q}, 0_{l}, 0_{l}\right)
$$

Let $m=l+s$ and $n=s+2 q+2 l$. Using the previous variables, (4.2) can be equivalently written as:

$$
\begin{align*}
\min _{x} & \frac{1}{2} x^{\top} Q x+c^{\top} x \\
\text { s.t. } & A x=b  \tag{4.3}\\
& x_{\mathcal{I}} \geq 0, x_{\mathcal{F}} \text { free, } \mathcal{I}=\{s+1, \ldots, n\}, \mathcal{F}=\{1, \ldots, s\}
\end{align*}
$$

where $b=0_{s+l} \in \mathbb{R}^{m}$,

$$
x=\left[u^{\top},\left(w^{+}\right)^{\top},\left(w^{-}\right)^{\top},\left(d^{+}\right)^{\top},\left(d^{-}\right)^{\top}\right]^{\top}, \quad c=\left[-\frac{\hat{y}^{\top}}{s}, \tau_{1} e_{w}^{\top}, \tau_{1} e_{w}^{\top}, \tau_{2} e_{d}^{\top}, \tau_{2} e_{d}^{\top}\right]^{\top} \in \mathbb{R}^{n}
$$

and

$$
Q=\left[\begin{array}{cc}
\frac{1}{s} I_{s} & 0_{s,(n-s)}  \tag{4.4}\\
0_{(n-s), s} & 0_{(n-s),(n-s)}
\end{array}\right] \in \mathbb{R}^{n \times n}, \quad A=\left[\begin{array}{ccccc}
-I_{s} & D & -D & 0_{s, l} & 0_{s, l} \\
0_{l, s} & L & -L & -I_{l} & I_{l}
\end{array}\right] \in \mathbb{R}^{m \times n}
$$

### 4.1 Specialized IP-PMM for Fused Lasso Least Squares

Notice that problem (4.3) is in the same form as (2.5). In what follows, we present a specialized inexact IP-PMM, suitable for solving unconstrained fused lasso least squares problems. The proposed specialized IP-PMM is characterized by the two following implementation details. Firstly, instead of factorizing system (2.8), we employ an iterative method (namely, the Preconditioned Conjugate Gradient (PCG)
method) to solve system (2.10). Secondly, as suggested in Section 3.1.1, we take advantage of the fact that the optimal solution of such problems is expected to be sparse, and use the heuristic approach that allows us to drop many of the variables of the problem, when the method is close to the optimal solution.

### 4.1.1 Solving the Newton System

We focus on solving the normal equations in (2.10). Let $k$ denote an arbitrary iteration of IP-PMM. We re-write the matrix in (2.10) without using the succinct notation introduced earlier:

$$
M_{k}=\left[\begin{array}{ll}
M_{1, k} & M_{2, k}^{\top}  \tag{4.5}\\
M_{2, k} & M_{3, k}
\end{array}\right]
$$

where:

$$
\begin{align*}
M_{1, k}= & \left(\left(\frac{1}{s}+\rho_{k}\right)^{-1}+\delta_{k}\right) I_{s}+D\left(\left(\Theta_{w^{+}, k}^{-1}+\rho_{k} I_{q}\right)^{-1}+\left(\Theta_{w^{-}, k}^{-1}+\rho_{k} I_{q}\right)^{-1}\right) D^{\top}, \\
M_{2, k}= & L\left(\Theta_{w^{+}}^{-1}+\rho_{k} I_{q}\right)^{-1} D^{\top}+L\left(\Theta_{w^{-}, k}^{-1}+\rho_{k} I_{q}\right)^{-1} D^{\top}, \\
M_{3, k}= & L\left(\left(\Theta_{w^{+}, k}^{-1}+\rho_{k} I_{q}\right)^{-1}+\left(\Theta_{w^{-}, k}+\rho_{k} I_{q}\right)^{-1}\right) L^{\top} \\
& +\left(\left(\Theta_{d^{+}, k}^{-1}+\rho_{k} I_{l}\right)^{-1}+\left(\Theta_{d^{-}, k}^{-1}+\rho_{k} I_{l}\right)^{-1}+\delta_{k} I_{l}\right), \tag{4.6}
\end{align*}
$$

and

$$
\left(\Theta_{k}^{\mathcal{T}}\right)^{-1}=\left[\begin{array}{cccc}
\Theta_{w^{+}, k}^{-1} & 0_{q, q} & 0_{q, l} & 0_{q, l} \\
0_{q, q} & \Theta_{w^{-}, k}^{-1} & 0_{q, l} & 0_{q, l} \\
0_{l, q} & 0_{l, q} & \Theta_{d^{+}, k}^{-1} & 0_{l, l} \\
0_{l, q} & 0_{l, q} & 0_{l, l} & \Theta_{d^{-}, k}^{-1}
\end{array}\right] .
$$

Notice that the matrix $D$ in (4.1) is dense, and hence we expect $M_{1, k}$ and $M_{2, k}$ in (4.6) to also be dense. On the other hand, $M_{3, k}$ remains sparse, and we know that $l \gg s$. As a consequence, the Cholesky factors of the matrix in (4.5) would inevitably contain dense blocks. Hence, it might be prohibitively expensive to compute such a decomposition. Instead, we solve the previous system using a PCG method. In order to do so efficiently, we must find an approximation for the coefficient matrix in (4.5). Given the fact that $M_{3, k}$ is sparse, while $M_{1, k}$ and $M_{2, k}$ are dense, we would like to find an approximation for the dense blocks. A possible approach would be to approximate $D$ by a low-rank matrix. Instead, based on the assumption $l \gg s$, we can approximate $M_{k}$ by the following block-diagonal preconditioner:

$$
P_{k}=\left[\begin{array}{cc}
M_{1, k} & 0_{s, l}  \tag{4.7}\\
0_{l, s} & M_{3, k}
\end{array}\right] \text {, where, } P_{k}^{-1}=\left[\begin{array}{cc}
M_{1, k}^{-1} & 0_{s, l} \\
0_{l, s} & M_{3, k}^{-1}
\end{array}\right] \text {. }
$$

Notice that $M_{3, k}$ does have a sparse Cholesky factor, due to the sparsity displayed in the discrete anisotropic TV matrix $L$. On the other hand, the Cholesky factors of $M_{1, k}$ are dense. However, computation and storage of these dense factors is possible, as we only need to perform $O\left(s^{3}\right)$ operations, and store $O\left(s^{2}\right)$ elements.

Spectral Analysis Let us now further support the choice of the preconditioner in (4.7) by performing a spectral analysis of the preconditioned system $R_{k}=P_{k}^{-1} M_{k}$. In the following, we write $\mathcal{A} \times \mathcal{B}$ to denote a vector space whose elements are vectors [ $\left.a^{\top}, b^{\top}\right]^{\top}$ with $a \in \mathcal{A}$ and $b \in \mathcal{B}$.

Theorem 4.1. Let $D \in \mathbb{R}^{s \times q}$ be the matrix in (4.1). Let also $M_{k} \in \mathbb{R}^{(s+l) \times(s+l)}$ be the matrix defined in (4.5) and $P_{k}$ the preconditioner defined in (4.7). Then, the preconditioned matrix $R_{k}=P_{k}^{-1} M_{k}$ has $l-\operatorname{rank}(D)$ eigenvalues $\lambda=1$, whose respective eigenvectors form a basis for $\left\{0_{s}\right\} \times\left\{\operatorname{Null}\left(M_{2, k}^{\top}\right)\right\}$. All the remaining eigenvalues of the preconditioned matrix satisfy $\lambda \in(\chi, 1) \cup(1,2)$, where $\chi=\frac{\delta_{k} \rho_{k}}{\sigma_{\max }^{2}(A)+\rho_{k} \delta_{k}}, \delta_{k}, \rho_{k}$ are the regularization parameters of IP-PMM and $A$ is defined in (4.4).

Proof. Let us consider the following generalized eigenproblem:

$$
M_{k} p=\lambda P_{k} p .
$$

We partition the eigenvector $p$ as $p=\left[p_{s}^{\top}, p_{l}^{\top}\right]^{\top}$. Using (4.5), and (4.7), the eigenproblem can be written as:

$$
\begin{align*}
& p_{s}+M_{1, k}^{-1} M_{2, k}^{\top} p_{l}=\lambda p_{s} \\
& M_{3, k}^{-1} M_{2, k} p_{s}+p_{l}=\lambda p_{l} . \tag{4.8}
\end{align*}
$$

From (4.6) and (4.7) we know that $M_{k} \succ 0$ and $P_{k} \succ 0$ and hence $\lambda>0$. Let us separate the analysis in two cases.

Case 1: $\lambda=1$. This is the case for every $p$ such that

$$
p \in \mathcal{S}_{k}=\left\{0_{s}\right\} \times\left\{\operatorname{Null}\left(M_{2, k}^{\top}\right)\right\},
$$

which trivially satisfies (4.8) for $\lambda=1$. Upon noticing that

$$
\operatorname{dim}\left(\mathcal{S}_{k}\right)=\operatorname{dim}\left(\operatorname{Null}\left(M_{2, k}^{\top}\right)\right)=l-\operatorname{rank}\left(M_{2, k}^{\top}\right),
$$

we can conclude that $R_{k}=P_{k}^{-1} M_{k}$ has an eigenvalue $\lambda=1$ of multiplicity $l-$ $\operatorname{rank}\left(M_{2, k}^{\top}\right)$. The respective eigenvectors form a basis of $\mathcal{S}_{k}$. Notice also, from (4.6), that $\operatorname{rank}\left(M_{2, k}^{\top}\right)=\operatorname{rank}(D) \leq s$.

Case 2: $\lambda \neq 1$. There are exactly $s+\operatorname{rank}(D)$ such eigenvalues. In order to analyze this case, we have to consider two generalized eigenproblems. On the one hand, from the first block equation in (4.8), we have that:

$$
p_{s}=\frac{1}{\lambda-1} M_{1, k}^{-1} M_{2, k}^{\top} p_{l},
$$

and hence, substituting this in the second block equation of (4.8) gives the following generalized eigenproblem:

$$
\begin{equation*}
G_{l, k} p_{l}=\nu M_{3, k} p_{l}, \tag{4.9}
\end{equation*}
$$

where $G_{l, k}=M_{2, k} M_{1, k}^{-1} M_{2, k}^{\top}$ and $\nu=(\lambda-1)^{2}$. By assumption $\lambda \neq 1$, and hence $\lambda= \pm \sqrt{\nu}+1$. However, we have $M_{k} \succ 0$ and hence $M_{3, k}-M_{2, k} M_{1, k}^{-1} M_{2, k}^{\top} \succ 0$. Let us assume that the maximum eigenvalue of $M_{3, k}^{-1} G_{l, k}$ is greater than or equal to

1, i.e. $\nu_{\max } \geq 1$. By substituting $\nu_{\max }$ in (4.9) and multiplying both sides of the previous inequality by $p_{l}^{\top}$, we get

$$
p_{l}^{\top} G_{l, k} p_{l}=\nu_{\max } p_{l}^{\top} M_{3, k} p_{l}
$$

and hence

$$
p_{l}^{\top}\left(G_{l, k}-M_{3, k}\right) p_{l} \geq 0
$$

The latter contradicts the fact that $M_{3, k}-G_{l, k} \succ 0$, and thus $\nu_{\max }<1$. In other words, $\lambda \in(0,1) \cup(1,2)$.

Similarly, starting from the second block equation of (4.8), we get

$$
p_{l}=\frac{1}{\lambda-1} M_{3, k}^{-1} M_{2, k} p_{s}
$$

and substituting this in the first block equation in (4.8) yields

$$
G_{s, k} p_{s}=\nu M_{1, k} p_{s}
$$

where $G_{s, k}=M_{2, k}^{\top} M_{3, k}^{-1} M_{2, k}$ and $\nu=(\lambda-1)^{2}$. As before, we have that $\lambda= \pm \sqrt{\nu}+1$. Using the fact that $M_{1, k}-M_{2, k}^{\top} M_{3, k}^{-1} M_{2, k} \succ 0$, we can mirror the previous analysis to conclude that $\nu_{\max }<1$, and hence $\lambda \in(0,1) \cup(1,2)$.

Finally, notice that as long as the primal and dual regularization parameters of IP-PMM, $\rho_{k}$ and $\delta_{k}$ respectively, are bounded away from zero, so are the eigenvalues of $R_{k}=P_{k}^{-1} M_{k}$, independently of the iteration $k$ of the algorithm. In particular, we have that

$$
\lambda_{\min }\left(R_{k}\right) \geq \frac{\lambda_{\min }\left(M_{k}\right)}{\lambda_{\max }\left(P_{k}\right)} \geq \frac{\delta_{k} \rho_{k}}{\sigma_{\max }^{2}(A)+\rho_{k} \delta_{k}}=\chi
$$

where we used the fact that $\lambda_{\min }\left(M_{k}\right) \geq \delta_{k}$ and $\lambda_{\max }\left(P_{k}\right) \leq \frac{\sigma_{\max }^{2}(A)}{\rho_{k}}+\delta_{k}$, where $A$ is defined in (4.4).

### 4.1.2 Dropping Primal Variables

The preconditioner (4.7) may be computed (and applied) very efficiently as we expect the Cholesky factor of $M_{3, k}$ to preserve sparsity and $M_{1, k} \in \mathbb{R}^{s \times s}$ to be relatively small (recall that $s \ll l$ ). We deduce from Theorem 4.1 that the preconditioner defined in (4.7) remains effective as long as the regularization parameters $\rho_{k}$ and $\delta_{k}$ are not too small. However, to attain convergence of IP-PMM $\rho_{k}$ and $\delta_{k}$ have to be reduced and then, due to the nature of IPMs, the matrix in (2.10) becomes increasingly ill conditioned as the method approaches the optimal solution. This implies that the preconditioner defined in (4.7). has only a limited applicability. In particular, this means that there is a limited scope for refining it and we may not be able to prevent degrading behaviour of PCG when IPM gets very close to the optimal solution.

However, we notice that the optimal solution of problem (4.2) is expected to be sparse. Like in the portfolio optimization problem, in light of the reformulation (4.3), we know that most of the primal variables $x$ converge to zero. Close to optimality the presence of such variables would adversely affect the conditioning of the matrix in (2.10). To prevent that, we employ a heuristic similar to the one introduced in Section 3.1.1 which consists of eliminating variables which approach zero and have
an associated Lagrange multiplier bounded away from zero. Given $\epsilon_{\text {drop }}>0, \xi>0$ and $\mathcal{V}=\emptyset$, at every iteration $k$ of IP-PMM, we add to $\mathcal{V}$ each variable $j \in \mathcal{I} \backslash \mathcal{V}$ satisfying condition (3.5), and replace (2.10) with the reduced system

$$
\begin{align*}
\left(A ^ { \mathcal { H } , \mathcal { G } } \left(Q^{\mathcal{G}, \mathcal{G}}+\left(\Theta_{k}^{\mathcal{G}, \mathcal{G}}\right)^{-1}\right.\right. & \left.\left.+\rho_{k} I_{|\mathcal{G}|}\right)^{-1}\left(A^{\mathcal{H}, \mathcal{G}}\right)^{\top}+\delta_{k} I_{m}\right) \Delta y \\
& =r_{2, k}+A^{\mathcal{H}, \mathcal{G}}\left(Q^{\mathcal{G}, \mathcal{G}}+\left(\Theta_{k}^{\mathcal{G}, \mathcal{G}}\right)^{-1}+\rho_{k} I_{|\mathcal{G}|}\right)^{-1} r_{1, k}^{\mathcal{G}} \tag{4.10}
\end{align*}
$$

where $\mathcal{H}=\{1, \ldots, m\}, \mathcal{G}=\mathcal{F} \cup(\mathcal{I} \backslash \mathcal{V})$ (assuming $\left(\Theta_{k}^{\mathcal{F}, \mathcal{F}}\right)^{-1}=0_{s, s}$ ), and $r_{1, k}, r_{2, k}$ are defined in (2.9) (with constraint matrix $A^{\mathcal{H}, \mathcal{G}}$ ).

### 4.2 Computational Experience

We consider a dataset consisting of fMRI scans for 16 male healthy US college students (age 20 to 25 ), with the aim of analyzing two active conditions: viewing unpleasant and pleasant images [44]. The preprocessed and registered data, available from https://github.com/lucabaldassarre/neurosparse, consist of 1344 scans of size 122,128 voxels (only voxels with probability greater than 0.5 of being in the gray matter are considered), with 42 scans considered per subject and active condition (i.e., 84 scans per subject in total).

In order to assess the performance of the IP-PMM on this type of problems, we carry out a comparison with two state-of-the-art algorithms for the solution of problem (4.2):

- FISTA. As done for the tests in [4], problem (4.2) is reformulated as

$$
\min _{w} \frac{1}{2 s}\|D w-\hat{y}\|^{2}+\|\hat{L} w\|_{1}
$$

where $\hat{L}=\left[\begin{array}{ll}\tau_{1} I_{q} & \tau_{2} L^{\top}\end{array}\right]^{\top}$, and solved by a version of FISTA [6] in which the proximal operator associated with $\|\hat{L} w\|_{1}$ is approximated by 10 steps of an inner FISTA cycle.

- ADMM. We consider the ADMM method [11] applied to the problem

$$
\begin{aligned}
\min _{w, u, d} & \frac{1}{2 s}\|D w-\hat{y}\|^{2}+\tau_{1}\|u\|_{1}+\tau_{2}\|d\|_{1} \\
\text { s.t. } & w-u=0_{q} \\
& L w-d=0_{l}
\end{aligned}
$$

in which the minimization of the quadratic function associated with the update of $w$ is approximated by 10 steps of the CG algorithm.

In Table 3 we show the results obtained by applying the algorithms to the solution of the fMRI data classification problem. For each choice of the pair of regularization parameters $\left(\tau_{1}, \tau_{2}\right)$, we report the average results obtained in a Leave-One-SubjectOut (LOSO) cross-validation test over the full dataset of patients. Because of this setting, for each problem the size of $w$ is $q=122,128$, the number of rows of $D$ is $s=1260$, and the dimension of $d=L w$ is $l=339,553$.

By preliminary experiments the choice $\tau_{1}=\tau_{2}$ appeared the most appropriate. Furthermore, for the IP-PMM, the parameters $\epsilon_{\text {drop }}$ and $\xi$ controlling the heuristic described in Section 4.1 .2 are set to $10^{-6}$ and $10^{2}$, respectively. To perform a fair comparison between the three algorithms, we consider a stopping criterion based on the execution time, which, after some preliminary tests, is fixed to 30 minutes. The solution of the normal equations system (4.10) is computed by the MATLAB PCG function, for which we set the maximum number of iterations to 2000 and the tolerance as

$$
\text { tol }= \begin{cases}10^{-4} & \text { if }\left\|r_{y, k}\right\|<1 \\ \max \left\{10^{-8}, \frac{10^{-4}}{\left\|r_{y, k}\right\|}\right\} & \text { otherwise }\end{cases}
$$

where $r_{y, k}$ is the right-hand side of equation (4.10).
For each algorithm tested, we report the mean and the standard deviation for three quality measures of the solution: classification accuracy (ACC), solution density (DEN) and corrected pairwise overlap (CORR OVR) (see [4, Section 2.3.3]). The latter, which is maybe the less common in the field of machine learning, is meant to measure the "stability" of the voxel selection. The three metrics are computed after thresholding the solution, as in [4]: after sorting the entries by their increasing magnitude, we set to zero the entries contributing at most to $0.01 \%$ of the $\ell_{1}$-norm of the solution.

Table 3: Comparison of IP-PMM, FISTA and ADMM in terms of the LOSO cross-validation scores

| Algorithm | $\tau_{1}=\tau_{2}$ | ACC | DEN | CORR OVR |
| :--- | ---: | ---: | ---: | ---: |
| IP-PMM | $10^{-2}$ | $86.16 \pm 7.11$ | $20.56 \pm 6.63$ | $43.47 \pm 9.09$ |
|  | $5 \cdot 10^{-2}$ | $84.90 \pm 4.80$ | $3.77 \pm 0.84$ | $62.70 \pm 10.39$ |
|  | $10^{-1}$ | $82.29 \pm 6.22$ | $2.49 \pm 0.34$ | $82.60 \pm 9.24$ |
| FISTA | $10^{-2}$ | $86.90 \pm 5.01$ | $88.97 \pm 0.71$ | $5.43 \pm 0.43$ |
|  | $5 \cdot 10^{-2}$ | $84.15 \pm 5.92$ | $19.36 \pm 0.86$ | $65.50 \pm 2.68$ |
|  | $10^{-1}$ | $81.62 \pm 7.58$ | $5.14 \pm 0.44$ | $80.44 \pm 9.72$ |
| ADMM | $10^{-2}$ | $86.46 \pm 6.91$ | $98.70 \pm 0.03$ | $0.03 \pm 0.01$ |
|  | $5 \cdot 10^{-2}$ | $85.57 \pm 5.37$ | $97.97 \pm 0.05$ | $0.15 \pm 0.04$ |
|  | $10^{-1}$ | $82.07 \pm 6.51$ | $97.50 \pm 0.19$ | $0.26 \pm 0.13$ |

By looking at Table 3, one can see that IP-PMM appears to be generally better than the other algorithms in enforcing the structured sparsity of the solution, presenting a good level of sparsity and overlap. It is worth noting that, because of its definition, the corrected pairwise overlap tends to zero as the density goes towards $100 \%$. Hence, for ADMM, which seems to be unable to enforce sparsity in the solution, the overlap is close to zero. As suggested in [4], one can evaluate the results in terms of the distance of the pair (ACC, CORR OVR) from the pair $(100,100)$ (the smaller the distance, the better the results). For the tests reported in the table, we can see that the best scores are obtained by IP-PMM with regularization parameters $\tau_{1}=\tau_{2}=10^{-1}$, for which the average accuracy is $82.3 \%$ and the corrected overlap is $82.6 \%$ with an average solution density of $2.5 \%$.

To further evaluate the efficiency of IP-PMM in the solution of this class of problems, we compare its performance in terms of elapsed time against the performance


Figure 1: History of classification accuracy, solution density and corrected pairwise overlap for IP-PMM (left) and FISTA (right), in the case $\tau_{1}=\tau_{2}=10^{-1}$. For the three quantities we report average measures with $95 \%$ confidence intervals.
of FISTA on the problem where the two methods reach the best scores, i.e., with $\tau_{1}=\tau_{2}=10^{-1}$. For all the 16 instances of the LOSO cross validation, we store the current solution of each algorithm after every minute and, at the end of the execution, we compute the three quality measures for such intermediate solutions. The results are shown in Figure 1 in terms of history of the mean values (lines) together with their $95 \%$ confidence intervals (shaded regions). From the plots we can see that while FISTA reaches the measures reported in Table 3 at the end of the 30 -minute run, the performance of IP-PMM stabilizes after about 20 minutes.

## 5 TV-based Poisson Image Restoration

Next we consider the restoration of images corrupted by Poisson noise, which arises in many applications, such as fluorescence microscopy, computed tomography (CT) and astronomical imaging (see, e.g., [25] and the references therein). In the discrete formulation of the restoration problem, the object to be restored is represented by a vector $w \in \mathbb{R}^{n}$ and the measured data are assumed to be a vector $g \in \mathbb{N}_{0}^{m}$, whose entries $g^{j}$ are samples from $m$ independent Poisson random variables $G^{j}$ with probability

$$
P\left(G^{j}=g^{j}\right)=\frac{e^{-(D w+a)^{j}}\left[(D w+a)^{j}\right]^{g^{j}}}{g^{j}!}
$$

where $a \in \mathbb{R}_{+}^{m}$ models the background radiation detected by the sensors. The matrix $D=\left(d^{i j}\right) \in \mathbb{R}^{m \times n}$ models the functioning of the imaging system and satisfies

$$
d^{i j} \geq 0 \text { for all } i, j, \quad \sum_{i=1}^{m} d^{i j}=1 \text { for all } j
$$

Here we assume that $D$ represents a convolution operator with periodic boundary conditions, which implies that $D$ has a Block-Circulant structure with Circulant Blocks (BCCB). Hence, $D w$ is computed expeditiously using the 2-dimensional Fast Fourier Transform (FFT). The maximum-likelihood approach [9] for the estimation
of $u$ leads to the minimization of the Kullback-Leibler (KL) divergence of $D w+a$ from $g$ :

$$
\begin{equation*}
D_{K L}(w) \equiv D_{K L}(D w+a, g)=\sum_{j=1}^{m}\left(g^{j} \ln \frac{g^{j}}{(D w+a)^{j}}+(D w+a)^{j}-g^{j}\right) \tag{5.1}
\end{equation*}
$$

where we set $g^{j} \ln \left(g^{j} /(D w+a)^{j}\right)=0$ if $g^{j}=0$ (we implicitly assume that $g$ has been converted into a real vector with entries ranging in the same interval as the entries of $w$ ). Since the estimation problem is highly ill conditioned, a regularization term is added to (5.1). We consider the Total Variation (TV) [50], which has received considerable attention because of its ability of preserving edges and smoothing flat areas of the images. Notice that, while it may introduce staircase artifacts, TV is still applied in many medical and biological applications (see, e.g., [5, 43, 62] and http://ranger.uta.edu/~huang/R_CSMRI.htm). The feasible set of the problem is defined by non-negativity constraints on the image intensity and the linear constraint $\sum_{i=1}^{n} w^{i}=\sum_{j=1}^{m}\left(g^{j}-a^{j}\right) \equiv r$ which guarantees preservation of the total intensity of the image.

The resulting model is

$$
\begin{array}{cl}
\min _{w} & D_{K L}(w)+\lambda\|L w\|_{1} \\
\text { s.t. } & e_{n}^{\top} w=r  \tag{5.2}\\
& w \geq 0
\end{array}
$$

where $L \in \mathbb{R}^{l \times n}$ is the matrix arising from the discretization of the TV functional (as in [15]).

### 5.1 Specialized IP-PMM for Image Restoration Problems

By employing the splitting strategy used in the previous sections, we can transform problem (5.2) to the following equivalent form:

$$
\begin{array}{cl}
\min _{x} & f(x) \equiv D_{K L}(w)+c^{\top} u \\
\text { s.t. } & A x=b  \tag{5.3}\\
& x \geq 0
\end{array}
$$

where, after introducing the additional constraint $d=L w$, and letting $\bar{m}=l+1$, $\bar{n}=n+2 l$, we set $x=\left[w^{\top}, u^{\top}\right]^{\top} \in \mathbb{R}^{\bar{n}}, u=\left[\left(d^{+}\right)^{\top},\left(d^{-}\right)^{\top}\right]^{\top} \in \mathbb{R}^{2 l}, c=\lambda e_{2 l}$, $b=\left[\begin{array}{ll}r, & 0_{l}^{\top}\end{array}\right]^{\top} \in \mathbb{R}^{\bar{m}}$, and

$$
A=\left[\begin{array}{ccc}
e_{n}^{\top} & 0_{l}^{\top} & 0_{l}^{\top} \\
L & -I_{l} & I_{l}
\end{array}\right] \in \mathbb{R}^{\bar{m} \times \bar{n}}
$$

We solve problem (5.3) by using IP-PMM combined with a perturbed composite Newton method [53]. Following the presentation in Section 2.2, we know that at the $k$-th iteration of the method we have to solve two linear systems of the form of (2.8). In order to avoid factorizations, every such system is solved using the preconditioned MINimum RESidual (MINRES) method. In order to accelerate the convergence of MINRES, we employ a block-diagonal preconditioner, which uses a
diagonal approximation of $\nabla^{2} f(x)$. More specifically, at iteration $k$ of IP-PMM, we have the following coefficient matrix:

$$
M_{k}=\left[\begin{array}{cc}
-H_{k} & A^{\top} \\
A & \delta_{k} I_{\bar{m}}
\end{array}\right]
$$

where $H_{k}=\left(\nabla^{2} f\left(x_{k}\right)+\Theta_{k}^{-1}+\rho_{k} I_{\bar{n}}\right)$, and we precondition it using the matrix

$$
\widetilde{M}_{k}=\left[\begin{array}{cc}
\widetilde{H}_{k} & 0_{\bar{n}, \bar{m}}  \tag{5.4}\\
0_{\bar{m}, \bar{n}} & A \widetilde{H}_{k}^{-1} A^{\top}+\delta_{k} I_{\bar{m}}
\end{array}\right]
$$

where $\widetilde{H}_{k}$ is a diagonal approximation of $H_{k}$. In order to analyze the spectral properties of the preconditioned matrix, we follow the developments in [8]. More specifically, we define $\widehat{H}_{k}:=\widetilde{H}_{k}^{-\frac{1}{2}} H_{k} \widetilde{H}_{k}^{\frac{1}{2}}$, and let:

$$
\alpha_{H}=\lambda_{\min }\left(\widehat{H}_{k}\right), \quad \beta_{H}=\lambda_{\max }\left(\widehat{H}_{k}\right), \quad \kappa_{H}=\frac{\beta_{H}}{\alpha_{H}}
$$

Using this notation, we know that an arbitrary element of the numerical range of this matrix is represented as $\gamma_{H} \in W\left(\widehat{H}_{k}\right)=\left[\alpha_{H}, \beta_{H}\right]$. Furthermore, we observe that in the special case where $\widetilde{H}_{k}=\operatorname{diag}\left(H_{k}\right)$, we have $\alpha_{H} \leq 1 \leq \beta_{H}$ since

$$
\frac{1}{n+2 l} \sum_{i=1}^{n+2 l} \lambda_{i}\left(\widehat{H}_{k}^{-1} H_{k, j}\right)=\frac{1}{n+2 l} \operatorname{Tr}\left(\hat{H}_{k}^{-1} H_{k}\right)=1
$$

Theorem 5.1. Let $k$ be an arbitrary IP-PMM iteration. Then, the eigenvalues of $\widetilde{M}_{k}^{-1} M_{k}$ lie in the union of the following intervals:

$$
I_{-}=\left[-\beta_{H}-1,-\alpha_{H}\right], \quad I_{+}=\left[\frac{1}{1+\beta_{H}}, 1\right]
$$

Proof. The proof follows exactly the developments in [8, Theorem 3.3].
In problem (5.3), $f(x)=D_{K L}(w)+c^{\top} u$ and hence

$$
\nabla f(x)=\left[\begin{array}{c}
\nabla D_{K L}(w) \\
c
\end{array}\right], \quad \nabla^{2} f(x)=\left[\begin{array}{cc}
\nabla^{2} D_{K L}(w) & 0_{n, 2 l} \\
0_{2 l, n} & 0_{2 l, 2 l}
\end{array}\right]
$$

where

$$
\nabla D_{K L}(w)=D^{\top}\left(e_{m}-\frac{g}{D w+a}\right), \quad \nabla^{2} D_{K L}(w)=D^{\top} U(w)^{2} D
$$

with $U(w)=\operatorname{diag}\left(\frac{\sqrt{g}}{D w+a}\right)$. Here the square root and the ratio are assumed to be component-wise. Notice that $D$ might be dense; however, as previously noted, its action can be computed via the FFT. Unfortunately, $D^{\top} U(w)^{2} D$ is not expected to be close to multilevel circulant. Even if it could be well-approximated by a multilevel circulant matrix, the scaling matrix of IP-PMM would destroy this structure. In other words, we use the structure of $D$ only when applying it to a vector. As a result, we only store the first column of $D$ and we use the FFT to apply this matrix to a vector. This allows us to compute the action of the Hessian easily.
Remark 1. The obvious choice would be to employ the approximation $\widetilde{H}_{k}=\operatorname{diag}\left(H_{k}\right)$, but the structure of the problem makes this choice rather expensive. A more efficient alternative is to use $\widetilde{H}_{k}=U\left(w_{k}\right)^{2}$, which is easier to compute and, as we will see in the following section, leads to good reconstruction results in practice.

### 5.2 Computational Experience

To evaluate the performance of the IP-PMM on this class of problems, we consider a set of three $256 \times 256$ grayscale images, which are presented in Figure 2. For each of


Figure 2: The three $256 \times 256$ grayscale images of the image restoration tests.
the three images we set up three restoration tests, where the images are corrupted by Poisson noise and $D$ represents one of the following blurs: Gaussian blur (GB), motion blur (MB), and out-of-focus blur (OF).

We compare the proposed method with the state-of-the-art Primal-Dual Algorithm with Linesearch (PDAL) method proposed in [39]. By following the example of [59, Algorithm 2], problem (5.2) is reformulated as

$$
\begin{equation*}
\min _{w} \max _{p, y} g^{\top} \ln (1+y)-y^{\top}(D w+a)-\lambda w^{\top} L^{\top} p+\chi_{\infty}(p)+\chi_{\mathcal{F}}(w) \tag{5.5}
\end{equation*}
$$

where $\chi_{\infty}$ is the characteristic function of the $\infty$-norm unit ball and $\chi_{\mathcal{F}}$ the characteristic function of the feasible set $\mathcal{F}$ of problem (5.2). It is worth noting that the PDAL algorithm for the solution of problem (5.5) requires at each step a projection on the feasible set $\mathcal{F}$, which is performed here by using the secant algorithm proposed by Dai and Fletcher in [22]. Concerning the parameters of PDAL, we use the same notation and tuning as in [39]. Following Section 6 of that paper, we set $\mu=0.7$, $\delta=0.99$ and $\beta=25$. The initial steplength is $\tau=\sqrt{1 / \omega}$, where $\omega$ is an estimate of $\left\|M^{\top} M\right\|$ and $M=\left[D^{\top} L^{\top}\right]^{\top}$ is the matrix linking the primal and dual variables. In the IP-PMM, we use the MINRES code by Michael Saunders and co-workers, available from https://web.stanford.edu/group/SOL/software/minres/, for which we set the relative tolerance $t o l=10^{-4}$ and the maximum number of iterations at each call equal to 20 . The regularization parameter $\lambda$ is determined by trial and error to minimize the Root Mean Square Error (RMSE) obtained by IP-PMM. For all the problems, the starting point is chosen to be the noisy and blurry image, i.e., $g$.

For all 9 tests we run 20 iterations of the IP-PMM method and let PDAL run for the same amount of time. In Figure 3 we report a comparison between the two algorithms in terms of elapsed time versus Root Mean Square Error (RMSE) in the solution of the 9 instances described above. As can be seen from the plots, the IP-PMM clearly outperforms PDAL on the instances with GB and OF (columns 1 and 3 , respectively, of Figure 3 ), while on the instances characterized by MB the two algorithms perform comparably.


Figure 3: Comparison between IP-PMM and PDAL in terms of Root Mean Square Error (RMSE) vs execution time in the solution of the 9 image restoration problems. From top to bottom, the rows refer to the cameraman instances, the house instances and the peppers instances, respectively. From left to right, the columns refer to the GB, MB and OF, respectively.

Table 4: Comparison between IP-PMM and PDAL in terms of RMSE, PSNR and MSSIM computed at the solutions provided by the two algorithms.

|  | IP-PMM |  |  | PDAL |  |  |
| :--- | :---: | :---: | :---: | :---: | :---: | :---: |
| Problem | RMSE | PSNR | MSSIM | RMSE | PSNR | MSSIM |
| cameraman - GB | $4.85 \mathrm{e}-02$ | $2.63 \mathrm{e}+01$ | $8.33 \mathrm{e}-01$ | $5.02 \mathrm{e}-02$ | $2.60 \mathrm{e}+01$ | $8.22 \mathrm{e}-01$ |
| cameraman - MB | $5.52 \mathrm{e}-02$ | $2.52 \mathrm{e}+01$ | $8.11 \mathrm{e}-01$ | $5.59 \mathrm{e}-02$ | $2.51 \mathrm{e}+01$ | $7.77 \mathrm{e}-01$ |
| cameraman - OF | $5.14 \mathrm{e}-02$ | $2.58 \mathrm{e}+01$ | $7.98 \mathrm{e}-01$ | $5.26 \mathrm{e}-02$ | $2.56 \mathrm{e}+01$ | $7.62 \mathrm{e}-01$ |
| house - GB | $9.71 \mathrm{e}-02$ | $2.03 \mathrm{e}+01$ | $7.51 \mathrm{e}-01$ | $9.88 \mathrm{e}-02$ | $2.01 \mathrm{e}+01$ | $6.92 \mathrm{e}-01$ |
| house - MB | $2.70 \mathrm{e}-02$ | $3.14 \mathrm{e}+01$ | $8.67 \mathrm{e}-01$ | $2.77 \mathrm{e}-02$ | $3.11 \mathrm{e}+01$ | $8.43 \mathrm{e}-01$ |
| house - OF | $3.80 \mathrm{e}-02$ | $2.84 \mathrm{e}+01$ | $8.33 \mathrm{e}-01$ | $4.09 \mathrm{e}-02$ | $2.78 \mathrm{e}+01$ | $7.70 \mathrm{e}-01$ |
| peppers - GB | $1.23 \mathrm{e}-01$ | $1.82 \mathrm{e}+01$ | $7.46 \mathrm{e}-01$ | $1.25 \mathrm{e}-01$ | $1.81 \mathrm{e}+01$ | $6.57 \mathrm{e}-01$ |
| peppers - MB | $8.76 \mathrm{e}-02$ | $2.12 \mathrm{e}+01$ | $8.90 \mathrm{e}-01$ | $8.78 \mathrm{e}-02$ | $2.11 \mathrm{e}+01$ | $8.72 \mathrm{e}-01$ |
| peppers - OF | $9.47 \mathrm{e}-02$ | $2.05 \mathrm{e}+01$ | $8.01 \mathrm{e}-01$ | $9.70 \mathrm{e}-02$ | $2.03 \mathrm{e}+01$ | $6.60 \mathrm{e}-01$ |

To better analyze the difference between the solutions provided by the two algorithms, one can look at Table 4, where we report the value of three scores: RMSE, Peak Signal-to-Noise Ratio (PSNR), and Mean Structural SIMilarity (MSSIM) [25]. It is worth noting that for RMSE smaller values are better, while for PSNR and MSSIM, higher values indicate better noise removal and perceived similarity between the restored and original image, respectively. From the table it is clear that in all the considered cases IP-PMM is able to produce a better restored image than PDAL, having always a larger MSSIM, also when the RMSE and PSNR values are comparable.

For the sake of space, we now restrict the comparison to the cases where the two algorithms seem to have reached equivalent solutions in terms of RMSE, to understand the differences in the restored images. We focus on the three instances in which $D$ represents MB (second column of Figure 3). In Figure 4 we report the results for cameraman, house and peppers with MB. By looking at the images one can see that those reconstructed by IP-PMM appear to be smoother (look, for example, at the sky in cameraman and house), which somehow indicates that the IP-PMM is better than PDAL in enforcing the TV regularization. Observe that this "visual" difference is reflected by the higher values of MSSIM reported for IP-PMM in Table 4.

## 6 Linear Classification via Regularized Logistic Regression

Finally, we deal with the problem of training a linear binary classifier. Let us consider a matrix $D \in \mathbb{R}^{n \times s}$ whose rows $d_{i}^{\top}$, with $i \in\{1, \ldots, n\}$, represent the training points, and a vector of labels $g \in\{-1,1\}^{n}$. In other words, we have a training set with $n$ binary-labeled samples and $s$ features. According to the logistic model, the conditional probability of having the label $g^{i}$ given the point $d_{i}$ has the form

$$
p_{l o g}(w)_{i}=P\left(g^{i} \mid d_{i}\right)=\frac{1}{1+e^{-g^{i} d_{i}^{\top} w}}
$$

where $w \in \mathbb{R}^{s}$ is the vector of weights determining the unbiased linear model under consideration. By following the maximum-likelihood approach, the weight vector $w$ can be obtained by maximizing the log-likelihood function or, equivalently, by minimizing the logistic loss function, i.e., by solving

$$
\min _{w} \phi(w) \equiv \frac{1}{n} \sum_{i=1}^{n} \phi_{i}(w), \quad \phi_{i}(w)=\log \left(1+e^{-g^{i} d_{i}^{\top} w}\right)
$$

To cope with the inherent ill-posedness of the estimation process, a regularization term is usually added to the previous model. For large-scale instances, where the features tend to be redundant, an $\ell^{1}$-regularization term is usually introduced to enforce sparsity in the solution, thus embedding feature selection in the training process. This results in the well-studied $\ell^{1}$-regularized logistic regression model:

$$
\begin{equation*}
\min _{w} \phi(w)+\tau\|w\|_{1} \tag{6.1}
\end{equation*}
$$

where $\tau>0$.


Figure 4: Results on cameraman, house and peppers with MB: noisy and blurry images (left), images restored by IP-PMM (center), images restored by PDAL (right).

As done in the previous sections, we can replace the nonsmooth model (6.1) with an equivalent smooth convex programming problem, i.e.,

$$
\begin{gather*}
\min _{x} f(x) \equiv \phi(w)+c^{\top} u, \\
\text { s.t. } A x=b,  \tag{6.2}\\
\\
u \geq 0,
\end{gather*}
$$

where, after introducing the additional constraint $u=w$, with $u=\left[\left(d^{+}\right)^{\top},\left(d^{-}\right)^{\top}\right]^{\top} \in$ $\mathbb{R}^{2 s}$, and letting $\bar{m}=s, \bar{n}=3 s$, we set $x=\left[w^{\top}, u^{\top}\right]^{\top} \in \mathbb{R}^{n}, c=\tau e_{2 s}, b=0_{\bar{m}}$, and $A \in \mathbb{R}^{\bar{m} \times \bar{n}}$ defined as $A=\left[\begin{array}{lll}I_{s} & -I_{s} & I_{s}\end{array}\right]$. The version of IP-PMM solving problem (6.2) is very similar to the one used to solve (5.3). The only difference here lies in the preconditioner. In particular, when solving problems of the form (6.2), we use the preconditioner defined in (5.4) (and subsequently analyzed in Theorem 5.1),
but we set $\widetilde{H}_{k}=\operatorname{diag}\left(H_{k}\right)$.

### 6.1 Computational Experience

To illustrate the performance of the IP-PMM on this class of problems, we consider a set of three linear classification problems from the LIBSVM dataset for binary classification (https://www.csie.ntu.edu.tw/~cjlin/libsvmtools/datasets/binary. $\mathrm{html})$. The names of the datasets, together with their number of features, training points and testing points are summarized in Table 5. For real-sim there is no predetermined separation of data between train and test, hence we apply a hold-out strategy keeping $30 \%$ of the data for testing.

Table 5: Characteristics of the $\ell^{1}$-regularized logistic regression problems

| Problem | Features | Train pts | Test pts |
| :--- | ---: | ---: | ---: |
| gisette | 5000 | 6000 | 1000 |
| rcv1 | 47,236 | 20,242 | 677,399 |
| real-sim | 20,958 | 50,617 | 21,692 |

To overcome the absence of the hyperplane bias in model (6.1), we add to the data matrices a further column with all ones, hence the resulting size of the problems is equal to $s+1$. For all the problems we set $\tau=\frac{1}{n}$, which is a standard choice in the literature.

To assess the effectiveness and efficiency of the proposed method we compare it with two state-of-the-art methods:

- an ADMM [11], available from http://www.stanford.edu/~boyd/papers/ distr_opt_stat_learning_admm.html;
- a MATLAB implementation of the newGLMNET method [60] used in LIBSVM, developed by the authors of [61] and available from https://github. com/ZiruiZhou/IRPN.

As in the tests presented in Section 5.2, the solution of the augmented system in IPPMM is performed by means of the MINRES implementation by Michael Saunders' team, with maximum number of iterations equal to 20 and tolerance $t o l=10^{-4}$.

We compare the three algorithms in terms of objective function value and classification error versus execution time, on runs lasting 15 seconds. The plots are reported in Figure 5. The IP-PMM is comparable with newGLMNET on the gisette instance, characterized by a very dense ( $>99 \%$ ) training data matrix, and both IP-PMM and newGLMNET clearly outperform ADMM. On the rcv1 and real-sim instances the IP-PMM method sightly outperforms newGLMNET in terms of classification error, and it is noticeably better in terms of the objective function value.

## 7 Conclusions

We have presented specialized IPMs for quadratic and general convex nonlinear optimization problems that model various sparse approximation instances. We have shown that by a proper choice of linear algebra solvers, which are a key issue in IPMs,


Figure 5: Results on the three $\ell^{1}$-regularized logistic regression problems (objective function value and classification error versus execution time).
we are able to efficiently solve the larger but smooth optimization problems coming from a standard reformulation of the original ones. This confirms the ability of IPMs to handle large sets of linear and non-negativity constraints. Computational experiments have been performed on diverse applications: multi-period portfolio selection, classification of fMRI data, restoration of blurry images corrupted by Poisson noise, and linear binary classification via regularized logistic regression. Comparisons with state-of-the-art first-order methods, which are widely used to tackle sparse approximation problems, have provided evidence that the presented IPM approach can offer
a noticeable advantage over those methods, especially when dealing with not-so-well conditioned problems.

We also believe that the results presented in this work may provide a basis for an in-depth analysis of the application of IPMs to many sparse approximation methods, and we plan to work in that direction in the future.

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