# Parallel Multi-Block ADMM with $o(1 / k)$ Convergence 

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$$
\begin{aligned}
& \text { Abstract This paper introduces a parallel and distributed algorithm for solving the following } \\
& \text { minimization problem with linear constraints: } \\
& \qquad \begin{array}{l}
\text { minimize } \\
f_{1}\left(\mathbf{x}_{1}\right)+\cdots+f_{N}\left(\mathbf{x}_{N}\right) \\
\text { subject to } A_{1} \mathbf{x}_{1}+\cdots+A_{N} \mathbf{x}_{N}=c, \\
\mathbf{x}_{1} \in \mathcal{X}_{1}, \ldots, \mathbf{x}_{N} \in \mathcal{X}_{N},
\end{array}
\end{aligned}
$$

where $N \geq 2, f_{i}$ are convex functions, $A_{i}$ are matrices, and $\mathcal{X}_{i}$ are feasible sets for variable $\mathbf{x}_{i}$. Our algorithm extends the alternating direction method of multipliers (ADMM) and decomposes the original problem into $N$ smaller subproblems and solves them in parallel at each iteration. This paper shows that the classic ADMM can be extended to the $N$-block Jacobi fashion and preserve convergence in the following two cases: (i) matrices $A_{i}$ are mutually near-orthogonal and have full column-rank, or (ii) proximal terms are added to the $N$ subproblems (but without any assumption on matrices $A_{i}$ ). In the latter case, certain proximal terms can let the subproblem be solved in more flexible and efficient ways. We show that $\left\|\mathbf{x}^{k+1}-\mathbf{x}^{k}\right\|_{M}^{2}$ converges at a rate of $o(1 / k)$ where $M$ is a symmetric positive semi-definte matrix. Since the parameters used in the convergence analysis are conservative, we introduce a strategy for automatically tuning the parameters to substantially accelerate

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our algorithm in practice. We implemented our algorithm (for the case ii above) on Amazon EC 2 and tested it on basis pursuit problems with $>300 \mathrm{~GB}$ of distributed data. This is the first time that successfully solving a compressive sensing problem of such a large scale is reported.

Keywords Alternating direction method of multipliers • ADMM • Parallel and distributed computing - Convergence rate

## 1 Introduction

Let $N \geq 2$ be an integer. We consider the following convex optimization problem with $N$ blocks of variables:

$$
\begin{equation*}
\min _{\mathbf{x}_{1}, \mathbf{x}_{2}, \ldots, \mathbf{x}_{N}} \sum_{i=1}^{N} f_{i}\left(\mathbf{x}_{i}\right) \quad \text { s.t. } \sum_{i=1}^{N} A_{i} \mathbf{x}_{i}=c, \tag{1.1}
\end{equation*}
$$

where $\mathbf{x}_{i} \in \mathbb{R}^{n_{i}}, A_{i} \in \mathbb{R}^{m \times n_{i}}, c \in \mathbb{R}^{m}$, and $f_{i}: \mathbb{R}^{n_{i}} \rightarrow(-\infty,+\infty]$ are closed proper convex functions, $i=1,2, \ldots, N$. If an individual block is subject to constraint $\mathbf{x}_{i} \in \mathcal{X}_{i}$, where $\mathcal{X}_{i} \subseteq \mathbb{R}^{n_{i}}$ is a nonempty closed convex set, it can be incorporated in the objective function $f_{i}$ using the indicator function:

$$
I_{\mathcal{X}_{i}}\left(\mathbf{x}_{i}\right)=\left\{\begin{array}{cc}
0 & \text { if } \mathbf{x}_{i} \in \mathcal{X}_{i}  \tag{1.2}\\
+\infty & \text { otherwise }
\end{array}\right.
$$

In this study, we do not impose any additional assumptions such as strict convexity or differentiability on the objective functions $f_{i}$.

Such optimization problems arise from a broad spectrum of applications including signal and image processing, compressive sensing, statistics and machine learning. See [1-4,15, $16,29,31,35]$ for example.

In this paper, we focus on parallel and distributed optimization algorithms-Jacobi ADMM (Algorithm 3) and Jacobi-Proximal ADMM (Algorithm 4) below - for solving (1.1). Since both of the objective function and constraints of (1.1) are sums of terms on individual $\mathbf{x}_{i}$ 's (they are separable), the problem can be decomposed into $N$ smaller subproblems, which is solved in a parallel and distributed manner.

### 1.1 Literature Review

A simple distributed algorithm for solving (1.1) is dual decomposition [13], which is essentially a dual ascent method or dual subgradient method [34] as follows. Consider the Lagrangian for problem (1.1):

$$
\begin{equation*}
\mathcal{L}\left(\mathbf{x}_{1}, \ldots, \mathbf{x}_{N}, \lambda\right)=\sum_{i=1}^{N} f_{i}\left(\mathbf{x}_{i}\right)-\lambda^{\top}\left(\sum_{i=1}^{N} A_{i} \mathbf{x}_{i}-c\right) \tag{1.3}
\end{equation*}
$$

where $\lambda \in \mathbb{R}^{m}$ is the Lagrange multiplier or the dual variable. The method of dual decomposition iterates as follows: for $k \geq 1$,

$$
\left\{\begin{array}{l}
\left(\mathbf{x}_{1}^{k+1}, \mathbf{x}_{2}^{k+1}, \ldots, \mathbf{x}_{N}^{k+1}\right)=\arg \min _{\left\{\mathbf{x}_{i}\right]} \mathcal{L}\left(\mathbf{x}_{1}, \ldots, \mathbf{x}_{N}, \lambda^{k}\right),  \tag{1.4}\\
\lambda^{k+1}=\lambda^{k}-\alpha_{k}\left(\sum_{i=1}^{N} A_{i} \mathbf{x}_{i}^{k+1}-c\right)
\end{array}\right.
$$

where $\alpha_{k}>0$ is a step-size. Since all the $\mathbf{x}_{i}$ 's are separable in the Lagrangian function (1.3), the $\mathbf{x}$-update step reduces to solving $N$ individual $\mathbf{x}_{i}$-subproblems:

$$
\begin{equation*}
\mathbf{x}_{i}^{k+1}=\underset{\mathbf{x}_{\mathbf{i}}}{\arg \min } f_{i}\left(\mathbf{x}_{i}\right)-\left\langle\lambda^{k}, A_{i} \mathbf{x}_{i}\right\rangle, \quad \text { for } i=1,2, \ldots, N, \tag{1.5}
\end{equation*}
$$

and thus they can be carried out in parallel. With suitable choice of $\alpha_{k}$ and certain assumptions, dual decomposition is guaranteed to converge to an optimal solution [34]. However, the convergence of such subgradient method often tends to be slow in practice. Its convergence rate for general convex problems is $O(1 / \sqrt{k})$. The dual smoothing method [30] can be applied under certain conditions and improve the rate to $O(1 / k)$.

Another effective distributed approach is based on the alternating direction method of multipliers (ADMM). ADMM was introduced in $[14,16]$ to solve the special case of problem (1.1) with two blocks of variables $(N=2)$. It utilizes the augmented Lagrangian for (1.1):

$$
\begin{equation*}
\mathcal{L}_{\rho}\left(\mathbf{x}_{1}, \ldots, \mathbf{x}_{N}, \lambda\right)=\sum_{i=1}^{N} f_{i}\left(\mathbf{x}_{i}\right)-\lambda^{\top}\left(\sum_{i=1}^{N} A_{i} \mathbf{x}_{i}-c\right)+\frac{\rho}{2}\left\|\sum_{i=1}^{N} A_{i} \mathbf{x}_{i}-c\right\|_{2}^{2} \tag{1.6}
\end{equation*}
$$

which incorporates a quadratic penalty of the constraints (with a parameter $\rho>0$ ) into the Lagrangian. In each iteration, the augmented Lagrangian is minimized over $\mathbf{x}_{1}$ and $\mathbf{x}_{2}$ separately, one after the other, followed by a dual update for $\lambda$. The iterative scheme of ADMM is outlined below:

$$
\left\{\begin{array}{l}
\mathbf{x}_{1}^{k+1}=\arg \min _{\mathbf{x}_{1}} \mathcal{L}_{\rho}\left(\mathbf{x}_{1}, \mathbf{x}_{2}^{k}, \lambda^{k}\right)  \tag{1.7}\\
\mathbf{x}_{2}^{k+1}=\arg \min _{\mathbf{x}_{2}} \mathcal{L}_{\rho}\left(\mathbf{x}_{1}^{k+1}, \mathbf{x}_{2}, \lambda^{k}\right), \\
\lambda^{k+1}=\lambda^{k}-\rho\left(A_{1} \mathbf{x}_{1}^{k+1}+A_{2} \mathbf{x}_{2}^{k+1}-c\right)
\end{array}\right.
$$

To solve the problem (1.1) with $N \geq 3$ using ADMM, one can first convert the multi-block problem into an equivalent two-block problem via variable splitting [2]:

$$
\begin{align*}
& \min _{\left\{\mathbf{x}_{i}\right\},\left\{\mathbf{z}_{i}\right\}} \sum_{i=1}^{N} f_{i}\left(\mathbf{x}_{i}\right)+I_{\mathcal{Z}}\left(\mathbf{z}_{1}, \ldots, \mathbf{z}_{N}\right) \\
& \text { s.t. } A_{i} \mathbf{x}_{i}-\mathbf{z}_{i}=\frac{c}{N}, \quad \forall i=1,2, \ldots, N, \tag{1.8}
\end{align*}
$$

where $I_{\mathcal{Z}}$ is a indicator function defined by (1.2), and the convex set $\mathcal{Z}$ is given by

$$
\mathcal{Z}=\left\{\left(\mathbf{z}_{1}, \ldots, \mathbf{z}_{N}\right): \sum_{i=1}^{N} \mathbf{z}_{i}=0\right\}
$$

The variables in (1.8) can be grouped into two blocks: $\mathbf{x}:=\left(\mathbf{x}_{1}, \ldots, \mathbf{x}_{N}\right)$ and $\mathbf{z}:=$ $\left(\mathbf{z}_{1}, \ldots, \mathbf{z}_{N}\right)$, so that ADMM can directly apply. The augmented Lagrangian for (1.8) is given by
$\mathcal{L}_{\rho}(\mathbf{x}, \mathbf{z}, \lambda)=\sum_{i=1}^{N} f_{i}\left(\mathbf{x}_{i}\right)+I_{\mathcal{Z}}(\mathbf{z})-\sum_{i=1}^{N} \lambda_{i}^{\top}\left(A_{i} \mathbf{x}_{i}-\mathbf{z}_{i}-\frac{c}{N}\right)+\frac{\rho}{2} \sum_{i=1}^{N}\left\|A_{i} \mathbf{x}_{i}-\mathbf{z}_{i}-\frac{c}{N}\right\|_{2}^{2}$.
Since all the $\mathbf{x}_{i}$ 's are now fully decoupled, the resulting $\mathbf{x}$-subproblem decomposes into $N$ individual $\mathbf{x}_{i}$-subproblems, which can be carried out in parallel.

```
Algorithm 1: Variable Splitting ADMM (VSADMM)
    Initialize \(\mathbf{x}^{0}, \lambda^{0}, \rho>0\);
    for \(k=0,1, \ldots\) do
        Update \(\mathbf{z}_{i}\) then \(\mathbf{x}_{i}\) for \(i=1, \ldots, N\) in parallel by:
        \(\mathbf{z}_{i}^{k+1}=\left(A_{i} \mathbf{x}_{i}^{k}-\frac{c}{N}-\frac{\lambda_{i}^{k}}{\rho}\right)-\frac{1}{N}\left\{\sum_{j=1}^{N} A_{j} \mathbf{x}_{j}^{k}-\frac{c}{N}-\frac{\lambda_{j}^{k}}{\rho}\right\} ;\)
        \(\mathbf{x}_{i}^{k+1}=\arg \min _{\mathbf{x}_{i}} f_{i}\left(\mathbf{x}_{i}\right)+\frac{\rho}{2}\left\|A_{i} \mathbf{x}_{i}-\mathbf{z}_{i}^{k+1}-\frac{c}{N}-\frac{\lambda_{i}^{k}}{\rho}\right\|_{2}^{2} ;\)
        Update \(\lambda_{i}^{k+1}=\lambda_{i}^{k}-\rho\left(A_{i} \mathbf{x}_{i}^{k+1}-\mathbf{z}_{i}^{k+1}-\frac{c}{N}\right), \forall i=1, \ldots, N\).
```

The resulting $\mathbf{z}$-subproblem is a simple quadratic problem:

$$
\begin{equation*}
\mathbf{z}^{k+1}=\arg \min _{\left\{\mathbf{z}: \sum_{i=1}^{N} \mathbf{z}_{i}=0\right\}} \sum_{i=1}^{N} \frac{\rho}{2}\left\|A_{i} \mathbf{x}_{i}-\mathbf{z}_{i}-\frac{c}{N}-\frac{\lambda_{i}^{k}}{\rho}\right\|_{2}^{2} \tag{1.9}
\end{equation*}
$$

which admits a closed-form solution.
The distributed ADMM approach based on (1.8), by introducing splitting variables, substantially increases the number of variables and constraints in the problem, especially when $N$ is large. Thus, it is clear that this approach will not be very efficient. Another extension is to replace the two-block alternating minimization scheme by a sweep of Gauss-Seidel update, namely, update $\mathbf{x}_{i}$ for $i=1,2, \ldots, N$ sequentially as follows:

$$
\begin{align*}
\mathbf{x}_{i}^{k+1} & =\arg \min _{\mathbf{x}_{i}} \mathcal{L}_{\rho}\left(\mathbf{x}_{1}^{k+1}, \ldots, \mathbf{x}_{i-1}^{k+1}, \mathbf{x}_{i}, \mathbf{x}_{i+1}^{k}, \ldots, \mathbf{x}_{N}^{k}, \lambda^{k}\right) \\
& =\arg \min _{\mathbf{x}_{i}} f_{i}\left(\mathbf{x}_{i}\right)+\frac{\rho}{2}\left\|\sum_{j<i} A_{j} \mathbf{x}_{j}^{k+1}+A_{i} \mathbf{x}_{i}+\sum_{j>i} A_{j} \mathbf{x}_{j}^{k}-c-\frac{\lambda^{k}}{\rho}\right\|_{2}^{2} \tag{1.10}
\end{align*}
$$

```
Algorithm 2: Gauss-Seidel ADMM
    Initialize \(\mathbf{x}^{0}, \lambda^{0}, \rho>0\);
    for \(k=0,1, \ldots\) do
        Update \(\mathbf{x}_{i}\) for \(i=1, \ldots, N\) sequentially by:
        \(\mathbf{x}_{i}^{k+1}=\min _{\mathbf{x}_{i}} f_{i}\left(\mathbf{x}_{i}\right)+\frac{\rho}{2}\left\|\sum_{j<i} A_{j} \mathbf{x}_{j}^{k+1}+A_{i} \mathbf{x}_{i}+\sum_{j>i} A_{j} \mathbf{x}_{j}^{k}-c-\frac{\lambda^{k}}{\rho}\right\|_{2}^{2} ;\) Update
        \(\lambda^{k+1}=\lambda^{k}-\rho\left(\sum_{i=1}^{N} A_{i} \mathbf{x}_{i}^{k+1}-c\right)\).
```

Such Gauss-Seidel ADMM (Algorithm 2) has been considered lately, e.g., in [22,25]. However, it has been shown that the algorithm may not converge for $N \geq 3$ [5]. Although lack of convergence guarantee, some empirical studies show that Algorithm 2 is still very effective at solving many practical problems (see, e.g., [31,35,36]). Recent work has shown that additional assumptions (involving at least one strongly convex objective function) can guarantee convergence for $N=3$ with certain modifications to the algorithm $[6,11,18,26$, 27]. A disadvantage of Gauss-Seidel ADMM is that the blocks are updated one after another, which is not amenable for parallelization.

### 1.2 Jacobi-Type ADMM

The sequential nature of Gauss-Seidel ADMM motivates us to consider using a Jacobi-type scheme that updates all the $N$ blocks in parallel:

$$
\begin{align*}
\mathbf{x}_{i}^{k+1} & =\arg \min _{\mathbf{x}_{i}} \mathcal{L}_{\rho}\left(\mathbf{x}_{1}^{k}, \ldots, \mathbf{x}_{i-1}^{k}, \mathbf{x}_{i}, \mathbf{x}_{i+1}^{k}, \ldots, \mathbf{x}_{N}^{k}, \lambda^{k}\right) \\
& =\arg \min _{\mathbf{x}_{i}} f_{i}\left(\mathbf{x}_{i}\right)+\frac{\rho}{2}\left\|A_{i} \mathbf{x}_{i}+\sum_{j \neq i} A_{j} \mathbf{x}_{j}^{k}-c-\frac{\lambda^{k}}{\rho}\right\|_{2}^{2}, \forall i=1, \ldots, N \tag{1.11}
\end{align*}
$$

We present it in Algorithm 3 below.
The parallelization comes with a cost: this scheme is more likely to diverge than the Gauss-Seidel scheme when using the same parameter $\rho$. In fact, it may diverge even in the two-block case; see [21] for such an example. To guarantee its convergence, either additional assumptions or modifications to the algorithm must be made.

```
Algorithm 3: Jacobi ADMM
    Initialize \(\mathbf{x}^{0}, \lambda^{0}, \rho>0\);
    for \(k=0,1, \ldots\) do
        Update \(\mathbf{x}_{i}\) for \(i=1, \ldots, N\) in parallel by:
        \(\mathbf{x}_{i}^{k+1}=\arg \min _{\mathbf{x}_{i}} f_{i}\left(\mathbf{x}_{i}\right)+\frac{\rho}{2}\left\|A_{i} \mathbf{x}_{i}+\sum_{j \neq i} A_{j} \mathbf{x}_{j}^{k}-c-\frac{\lambda^{k}}{\rho}\right\|_{2}^{2} ;\)
        Update \(\lambda^{k+1}=\lambda^{k}-\rho\left(\sum_{i=1}^{N} A_{i} \mathbf{x}_{i}^{k+1}-c\right)\).
```

In Sect. 4, we show that if matrices $A_{i}$ are mutually near-orthogonal and have full columnrank, then Algorithm 3 converges globally. For general cases, a few variants of Jacobi ADMM with additional corrections were proposed in [20,21].

In this paper, we propose Jacobi-Proximal ADMM (Algorithm 4). Compared with Algorithm 3, we make no assumption to the problem data but add a proximal term $\frac{1}{2}\left\|\mathbf{x}_{i}-\mathbf{x}_{i}^{k}\right\|_{P_{i}}^{2}$ for each $\mathbf{x}_{i}$-subproblem and a damping parameter $\gamma>0$ for the update of $\lambda$. Here $P_{i} \succeq 0$ is a symmetric and positive semi-definite matrix and we let $\left\|\mathbf{x}_{i}\right\|_{P_{i}}^{2}:=\mathbf{x}_{i}^{\top} P_{i} \mathbf{x}_{i}$. We will show how to choose $P_{i}$ that ensure the convergence of Algorithm 4.

```
Algorithm 4: Jacobi-Proximal ADMM
    Initialize: \(\mathbf{x}_{i}^{0}(i=1,2, \ldots, N)\) and \(\lambda^{0}\);
    for \(k=0,1, \ldots\) do
        Update \(\mathbf{x}_{i}\) for \(i=1, \ldots, N\) in parallel by:
        \(\mathbf{x}_{i}^{k+1}=\arg \min _{\mathbf{x}_{i}} f_{i}\left(\mathbf{x}_{i}\right)+\frac{\rho}{2}\left\|A_{i} \mathbf{x}_{i}+\sum_{j \neq i} A_{j} \mathbf{x}_{j}^{k}-c-\frac{\lambda^{k}}{\rho}\right\|_{2}^{2}+\frac{1}{2}\left\|\mathbf{x}_{i}-\mathbf{x}_{i}^{k}\right\|_{P_{i}}^{2} ;\)
        Update \(\lambda^{k+1}=\lambda^{k}-\gamma \rho\left(\sum_{i=1}^{N} A_{i} \mathbf{x}_{i}^{k+1}-c\right)\).
```

The proposed algorithm has a few advantages. First of all, as we will show, it enjoys global convergence as well as an $o(1 / k)$ convergence rate of $\left\|\mathbf{x}^{k+1}-\mathbf{x}^{k}\right\|_{M}^{2}$ under certain conditions on $P_{i}$ and $\gamma$. Secondly, when the $\mathbf{x}_{i}$-subproblem is not strictly convex, adding the proximal term can make the subproblem strictly or strongly convex, making it to have a unique solution. Thirdly, there are multiple choices for matrices $P_{i}$ with which the subproblems become easier
to solve. Specifically, the $\mathbf{x}_{i}$-subproblem contains a quadratic term $\frac{\rho}{2} \mathbf{x}_{i}^{\top} A_{i}^{\top} A_{i} \mathbf{x}_{i}$. When $A_{i}^{\top} A_{i}$ is ill-conditioned or computationally expensive to invert, one can let $P_{i}=D_{i}-\rho A_{i}^{\top} A_{i}$, which cancels the quadratic term $\frac{\rho}{2} \mathbf{x}_{i}^{\top} A_{i}^{\top} A_{i} \mathbf{x}_{i}$ and adds $\frac{1}{2} \mathbf{x}_{i}^{\top} D_{i} \mathbf{x}_{i}$, where the matrix $D_{i}$ can be chosen as some well-conditioned and simple matrix (e.g., a diagonal matrix), thereby leading to an easier subproblem.

Let us mention two commonly used choices of $P_{i}$ :

- $P_{i}=\tau_{i} \mathbf{I}\left(\tau_{i}>0\right)$ : This corresponds to the standard proximal method.
- $P_{i}=\tau_{i} \mathbf{I}-\rho A_{i}^{\top} A_{i}\left(\tau_{i}>0\right)$ : This corresponds to the prox-linear method [7], which linearizes the quadratic penalty term of augmented Lagrangian at the current point $\mathbf{x}_{i}^{k}$ and adds a proximal term. The $\mathbf{x}_{i}$-subproblem is given by

$$
\begin{equation*}
\mathbf{x}_{i}^{k+1}=\arg \min _{\mathbf{x}_{i}} f_{i}\left(\mathbf{x}_{i}\right)+\left\langle\rho A_{i}^{\top}\left(A \mathbf{x}^{k}-c-\lambda^{k} / \rho\right), \mathbf{x}_{i}\right\rangle+\frac{\tau_{i}}{2}\left\|\mathbf{x}_{i}-\mathbf{x}_{i}^{k}\right\|^{2} \tag{1.12}
\end{equation*}
$$

It essentially uses an identity matrix $\tau_{i} \mathbf{I}$ to approximate the Hessian $\rho A_{i}^{\top} A_{i}$ of the quadratic term.
More choices of $P_{i}$ have also been discussed in [12,38].

### 1.3 Summary of Contributions

This paper introduces novel results from the following perspectives. Firstly, we propose Jacobi-Proximal ADMM (Algorithm 4), which is suitable for parallel and distributed computing.

The flexible use of proximal terms makes it possible to solve subproblems in different ways, important for easy coding and fast computation. We establish its convergence at a rate of $o(1 / k)$ for $\left\|\mathbf{x}^{k+1}-\mathbf{x}^{k}\right\|_{M}^{2}$. We also provide an adaptive parameter tuning scheme which updates the important parameters at little extra cost. Our numerical results on the exchange problem and $\ell_{1}$-minimization problem show that the algorithm achieves competitive performance.

The second contribution is a condition that guarantees the convergence of Jacobi ADMM (Algorithm 3). The condition is applied to the coefficient matrices $A_{i}$, without assumptions on the objective functions $f_{i}$ or penalty parameter $\rho$.

A minor contribution is the improvement of the established convergence rate of $O(1 / k)$ for the standard ADMM to $o(1 / k)$. The improvement requires showing a certain sequence is monotonically nonincreasing, which can be applied to improve the existing rate of $O(1 / k)$ of several other algorithms to $o(1 / k)$, e.g. [21].

Finally we have experimented Algorithm 4 numerically based on the Amazon Elastic Computing Cloud (EC2) to demonstrate its effectiveness of solving a compressive sensing problem of huge size.

### 1.4 Notation, Assumptions and Preliminary Results

To simplify the notation in this paper, we introduce

$$
\mathbf{x}:=\left(\begin{array}{c}
\mathbf{x}_{1} \\
\vdots \\
\mathbf{x}_{N}
\end{array}\right) \in \mathbb{R}^{\mathbf{n}}, \mathbf{A}:=\left(A_{1}, \ldots, A_{N}\right) \in \mathbb{R}^{\mathbf{m} \times \mathbf{n}}, \mathbf{u}:=\binom{\mathbf{x}}{\lambda} \in \mathbb{R}^{\mathbf{n}+\mathbf{m}}
$$

where $n=\sum_{i=1}^{N} n_{i}$. We let $\langle\cdot, \cdot\rangle$ and $\|\cdot\|$ denote the standard inner product and $\ell_{2}$-norm $\|\cdot\|_{2}$, respectively, in the Euclidean space. For a matrix $M \in \mathbb{R}^{l \times l},\|M\|$ denotes the spectral
norm, i.e., the largest singular value of $M$. For a positive definite matrix $G \in \mathbb{R}^{l \times l}$, we define the $G$-norm as follows:

$$
\begin{equation*}
\|\mathbf{z}\|_{G}:=\sqrt{\mathbf{z}^{\top} G \mathbf{z}}, \quad \forall \mathbf{z} \in \mathbb{R}^{l} . \tag{1.13}
\end{equation*}
$$

If the matrix $G$ is positive semi-definite, then $\|\cdot\|_{G}$ is a semi-norm.
Throughout the paper, we make the following standard assumptions.
Assumption 1 Functions $f_{i}: \mathbb{R}^{n_{i}} \rightarrow(-\infty,+\infty](i=1,2, \ldots, N)$ are closed proper convex.

Assumption 2 There exists a saddle point $\mathbf{u}^{*}=\left(\mathbf{x}_{1}^{*}, \mathbf{x}_{2}^{*}, \ldots, \mathbf{x}_{N}^{*}, \lambda^{*}\right)$ to the problem (1.1), that is, $\mathbf{u}^{*}$ satisfies the KKT conditions:

$$
\begin{align*}
& A_{i}^{\top} \lambda^{*} \in \partial f_{i}\left(\mathbf{x}_{i}^{*}\right), \text { for } i=1, \ldots, N  \tag{1.14a}\\
& A \mathbf{x}^{*}=\sum_{i=1}^{N} A_{i} \mathbf{x}_{i}^{*}=c \tag{1.14b}
\end{align*}
$$

The conditions (1.14a) and (1.14b) can be written in a more compact form using variational inequality [21]:

$$
\begin{equation*}
f(\mathbf{x})-f\left(\mathbf{x}^{*}\right)+\left(\mathbf{u}-\mathbf{u}^{*}\right)^{\top} F\left(\mathbf{u}^{*}\right) \geq 0, \forall \mathbf{u}, \tag{1.15}
\end{equation*}
$$

where $f(\mathbf{x}):=\sum_{i} f_{i}\left(\mathbf{x}_{i}\right)$ and $F(\mathbf{u}):=\left[-A_{1}^{\top} \lambda, \ldots,-A_{N}^{\top} \lambda, A \mathbf{x}-c\right]^{\top}$.
Let $\partial f_{i}\left(\mathbf{x}_{i}\right)$ denote the subdifferential of $f_{i}$ at $\mathbf{x}_{i}$ :

$$
\begin{equation*}
\partial f_{i}\left(\mathbf{x}_{i}\right):=\left\{s_{i} \in \mathbb{R}^{n_{i}}: s_{i}^{\top}\left(\mathbf{y}_{i}-\mathbf{x}_{i}\right) \leq f_{i}\left(\mathbf{y}_{i}\right)-f_{i}\left(\mathbf{x}_{i}\right), \forall \mathbf{y}_{i} \in \operatorname{dom} f_{i}\right\}, \tag{1.16}
\end{equation*}
$$

which is nonempty under Assumption 1. We recall that a subdifferential map of a convex function is monotone, which will be used several times in later sections.

Proposition 1.1 Under Assumption 1, for any $\mathbf{x}_{i}, \mathbf{y}_{i} \in \operatorname{dom} f_{i}$, we have

$$
\begin{equation*}
\left(s_{i}-t_{i}\right)^{\top}\left(\mathbf{x}_{i}-\mathbf{y}_{i}\right) \geq 0, \quad \forall s_{i} \in \partial f_{i}\left(\mathbf{x}_{i}\right), t_{i} \in \partial f_{i}\left(\mathbf{y}_{i}\right), i=1,2, \ldots, N . \tag{1.17}
\end{equation*}
$$

In addition, we use an elementary argument to improve the convergence rate from $O(1 / k)$ to $o(1 / k)$. Intuitively, the harmonic sequence $1 / k$ is not summable, so a summable, nonnegative, monotonic sequence shall converge faster than $1 / k$.

Lemma 1.1 If a sequence $\left\{a_{k}\right\} \subseteq \mathbb{R}$ obeys: (1) $a_{k} \geq 0$; (2) $\sum_{k=1}^{\infty} a_{k}<+\infty$; (3) $a_{k}$ is monotonically non-increasing, then we have $a_{k}=o(1 / k)$.

Proof Since $k \cdot a_{2 k} \leq a_{k+1}+a_{k+2}+\cdots+a_{2 k} \rightarrow 0$ as $k \rightarrow+\infty, a_{k}=o(1 / k)$.
More study on the $o(1 / k)$ rate, its tightness, and other results for splitting schemes including the standard ADMM can be found in [10].

### 1.5 Organization

The rest of the paper is organized as follows. In Sect. 2, we establish the convergence as well as an $o(1 / k)$ rate of convergence for the Jacobi-Proximal ADMM (Algorithm 4). We also propose a practical parameter tuning method, which updates important parameters on the fly at little extra cost and makes the algorithm run faster. In Sect. 3, we present several numerical results to demonstrate the efficiency of Algorithm 4 in comparison with some
existing parallel algorithms. In Sect. 4, we show that the Jacobi ADMM (Algorithm 3), though lack of convergence in the general case, will still converge under certain conditions on the matrices $A_{i}$. It is interesting to point out that the Jacobi-Proximal ADMM can be derived based on the classical proximal point method(PPA). See Sect. 5. Finally, we conclude the paper in Sect. 6. For convenience, we include some detail on how to use Lemma 1.1 to improve some existing convergence rates from $O(1 / k)$ to $o(1 / k)$ in "Appendix".

## 2 Convergence Analysis of the Jacobi-Proximal ADMM

In this section, we study the convergence of Jacobi-Proximal ADMM (Algorithm 4). We first show its convergence and establish an $o(1 / k)$ convergence rate for a quantity that is also used in [24]. Furthermore, we show how to automatically update the parameter in order to make Jacobi-Proximal ADMM more practical.

### 2.1 Convergence

To simplify the notation, we let

$$
G_{x}:=\left(\begin{array}{lll}
P_{1}+\rho A_{1}^{\top} A_{1} & & \\
& \ddots & \\
& & P_{N}+\rho A_{N}^{\top} A_{N}
\end{array}\right), G:=\left(\begin{array}{ll}
G_{x} & \\
& \frac{1}{\gamma \rho} \mathbf{I}
\end{array}\right),
$$

and

$$
Q:=\left(\begin{array}{cccc}
P_{1}+\rho A_{1}^{\top} A_{1} & & & \frac{1}{\gamma} A_{1}^{\top}  \tag{2.1}\\
& \ddots & \vdots \\
& & P_{N}+\rho A_{N}^{\top} A_{N} & \frac{1}{\gamma} A_{N}^{\top} \\
\frac{1}{\gamma} A_{1} & \cdots & \frac{1}{\gamma} A_{N} & \frac{2-\gamma}{\rho \gamma^{2}} \mathbf{I}
\end{array}\right),
$$

where $\mathbf{I}$ is the identity matrix of size $m \times m$. In the rest of the section, we let $\left\{\mathbf{u}^{k}=\right.$ $\left.\left(\mathbf{x}_{1}^{k}, \ldots, \mathbf{x}_{N}^{k}, \lambda^{k}\right), k \geq 1\right\}$ denote the sequence generated by Jacobi-Proximal ADMM from any initial point. The analysis is based on bounding the error $\left\|\mathbf{u}^{k}-\mathbf{u}^{*}\right\|_{G}^{2}$ and estimating its decrease, motivated by the works [12,21,23].

Lemma 2.1 For $k \geq 1$, we have

$$
\begin{equation*}
\left\|\mathbf{u}^{k}-\mathbf{u}^{*}\right\|_{G}^{2}-\left\|\mathbf{u}^{k+1}-\mathbf{u}^{*}\right\|_{G}^{2} \geq\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{Q}^{2}, \tag{2.2}
\end{equation*}
$$

where $\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{Q}^{2}:=\left\|\mathbf{x}^{k}-\mathbf{x}^{k+1}\right\|_{G_{x}}^{2}+\frac{2-\gamma}{\rho \gamma^{2}}\left\|\lambda^{k}-\lambda^{k+1}\right\|^{2}+\frac{2}{\gamma}\left(\lambda^{k}-\lambda^{k+1}\right)^{\top} A\left(\mathbf{x}^{k}-\mathbf{x}^{k+1}\right)$.
Proof Recall that in Algorithm 4, we solve the following $\mathbf{x}_{i}$-subproblem:

$$
\mathbf{x}_{i}^{k+1}=\arg \min _{\mathbf{x}_{i}} f_{i}\left(\mathbf{x}_{i}\right)+\frac{\rho}{2}\left\|A_{i} \mathbf{x}_{i}+\sum_{j \neq i} A_{j} \mathbf{x}_{j}^{k}-c-\frac{\lambda^{k}}{\rho}\right\|^{2}+\frac{1}{2}\left\|\mathbf{x}_{i}-\mathbf{x}_{i}^{k}\right\|_{P_{i}}^{2}
$$

Its optimality condition is given by

$$
\begin{equation*}
A_{i}^{\top}\left(\lambda^{k}-\rho\left(A_{i} \mathbf{x}_{i}^{k+1}+\sum_{j \neq i} A_{j} \mathbf{x}_{j}^{k}-c\right)\right)+P_{i}\left(\mathbf{x}_{i}^{k}-\mathbf{x}_{i}^{k+1}\right) \in \partial f_{i}\left(\mathbf{x}_{i}^{k+1}\right) \tag{2.3}
\end{equation*}
$$

For convenience, we let $\hat{\lambda}:=\lambda^{k}-\rho\left(A \mathbf{x}^{k+1}-c\right)$. Then (2.3) can be rewritten as

$$
\begin{equation*}
A_{i}^{\top}\left(\hat{\lambda}-\rho \sum_{j \neq i} A_{j}\left(\mathbf{x}_{j}^{k}-\mathbf{x}_{j}^{k+1}\right)\right)+P_{i}\left(\mathbf{x}_{i}^{k}-\mathbf{x}_{i}^{k+1}\right) \in \partial f_{i}\left(\mathbf{x}_{i}^{k+1}\right) . \tag{2.4}
\end{equation*}
$$

By Lemma 1.1, it follows from (1.14a) and (2.4) that

$$
\left\langle A_{i}\left(\mathbf{x}_{i}^{k+1}-\mathbf{x}_{i}^{*}\right), \hat{\lambda}-\lambda^{*}-\rho \sum_{j \neq i} A_{j}\left(\mathbf{x}_{j}^{k}-\mathbf{x}_{j}^{k+1}\right)\right\rangle+\left(\mathbf{x}_{i}^{k+1}-\mathbf{x}_{i}^{*}\right)^{\top} P_{i}\left(\mathbf{x}_{i}^{k}-\mathbf{x}_{i}^{k+1}\right) \geq 0 .
$$

Summing the above inequality over all $i$ and using the following equality for each $i$ :

$$
\sum_{j \neq i} A_{j}\left(\mathbf{x}_{j}^{k}-\mathbf{x}_{j}^{k+1}\right)=A\left(\mathbf{x}^{k}-\mathbf{x}^{k+1}\right)-A_{i}\left(\mathbf{x}_{i}^{k}-\mathbf{x}_{i}^{k+1}\right),
$$

we obtain

$$
\begin{align*}
& \left\langle A\left(\mathbf{x}^{k+1}-\mathbf{x}^{*}\right), \hat{\lambda}-\lambda^{*}\right\rangle+\sum_{i=1}^{N}\left(\mathbf{x}_{i}^{k+1}-\mathbf{x}_{i}^{*}\right)^{\top}\left(P_{i}+\rho A_{i}^{\top} A_{i}\right)\left(\mathbf{x}_{i}^{k}-\mathbf{x}_{i}^{k+1}\right) \\
& \quad \geq \rho\left\langle A\left(\mathbf{x}^{k+1}-\mathbf{x}^{*}\right), A\left(\mathbf{x}^{k}-\mathbf{x}^{k+1}\right)\right\rangle \tag{2.5}
\end{align*}
$$

Note that $A\left(\mathbf{x}^{k+1}-\mathbf{x}^{*}\right)=\frac{1}{\gamma \rho}\left(\lambda^{k}-\lambda^{k+1}\right)$ and

$$
\hat{\lambda}-\lambda^{*}=\left(\hat{\lambda}-\lambda^{k+1}\right)+\left(\lambda^{k+1}-\lambda^{*}\right)=\frac{\gamma-1}{\gamma}\left(\lambda^{k}-\lambda^{k+1}\right)+\left(\lambda^{k+1}-\lambda^{*}\right) .
$$

With the above two equations, the inequality (2.5) can be rewritten as

$$
\begin{align*}
& \left\langle\frac{1}{\gamma \rho}\left(\lambda^{k}-\lambda^{k+1}\right), \lambda^{k+1}-\lambda^{*}\right\rangle+\sum_{i=1}^{N}\left(\mathbf{x}_{i}^{k+1}-\mathbf{x}_{i}^{*}\right)^{\top}\left(P_{i}+\rho A_{i}^{\top} A_{i}\right)\left(\mathbf{x}_{i}^{k}-\mathbf{x}_{i}^{k+1}\right) \\
& \geq \frac{1-\gamma}{\gamma^{2} \rho}\left\|\lambda^{k}-\lambda^{k+1}\right\|^{2}+\frac{1}{\gamma}\left(\lambda^{k}-\lambda^{k+1}\right)^{\top} A\left(\mathbf{x}^{k}-\mathbf{x}^{k+1}\right) \tag{2.6}
\end{align*}
$$

or more compactly,

$$
\begin{align*}
\left(\mathbf{u}^{k}-\mathbf{u}^{k+1}\right)^{\top} G\left(\mathbf{u}^{k+1}-\mathbf{u}^{*}\right) \geq & \frac{1-\gamma}{\gamma^{2} \rho}\left\|\lambda^{k}-\lambda^{k+1}\right\|^{2} \\
& +\frac{1}{\gamma}\left(\lambda^{k}-\lambda^{k+1}\right)^{\top} A\left(\mathbf{x}^{k}-\mathbf{x}^{k+1}\right) . \tag{2.7}
\end{align*}
$$

Since $\left\|\mathbf{u}^{k}-\mathbf{u}^{*}\right\|_{G}^{2}-\left\|\mathbf{u}^{k+1}-\mathbf{u}^{*}\right\|_{G}^{2}=2\left(\mathbf{u}^{k}-\mathbf{u}^{k+1}\right)^{\top} G\left(\mathbf{u}^{k+1}-\mathbf{u}^{*}\right)+\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{G}^{2}$, using the above inequality (2.7) yields (2.2) immediately.

If the matrix $Q$ is positive definite, there exists some $\eta>0$ such that

$$
\begin{equation*}
\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{Q}^{2} \geq \eta \cdot\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|^{2} \geq 0 \tag{2.8}
\end{equation*}
$$

Then Lemma 2.1 indicates that

$$
\begin{equation*}
\left\|\mathbf{u}^{k}-\mathbf{u}^{*}\right\|_{G}^{2}-\left\|\mathbf{u}^{k+1}-\mathbf{u}^{*}\right\|_{G}^{2} \geq \eta \cdot\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|^{2} \tag{2.9}
\end{equation*}
$$

i.e., the iterative sequence $\left\{\mathbf{u}^{k}\right\}$ is strictly contractive. In particular, the error $\left\|\mathbf{u}^{k}-\mathbf{u}^{*}\right\|_{G}^{2}$ is monotonically non-increasing and thus converging, as well as $\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|^{2} \rightarrow 0$. Then the
convergence of the algorithm $\left(\left\|\mathbf{u}^{k}-\mathbf{u}^{*}\right\|_{G}^{2} \rightarrow 0\right)$ follows immediately from the standard analysis for contraction methods (see, e.g., [19]). We omit the details of the proof for the sake of brevity and state our convergence theorem as follows.

Theorem 2.1 Under Assumptions 1 and 2, one can find parameters in Algorithm 4 such that its sequence $\left\{\mathbf{u}^{k}\right\}$ converges to a solution $\mathbf{u}^{*}$ to the problem (1.1). Specifically, if one chooses some $\epsilon_{i}>0$ such that the parameters $\rho, \gamma$ and $P_{i}(i=1,2, \ldots, N)$ satisfy the following condition:

$$
\left\{\begin{array}{l}
P_{i} \succ \rho\left(\frac{1}{\epsilon_{i}}-1\right) A_{i}^{\top} A_{i}, i=1,2, \ldots, N  \tag{2.10}\\
\sum_{i=1}^{N} \epsilon_{i}<2-\gamma,
\end{array}\right.
$$

then $Q$ in (2.1) is positive definite and $\left\{\mathbf{u}^{k}\right\}$ converges to $\mathbf{u}^{*}$.
Furthermore, by letting each $\epsilon_{i}<\frac{2-\gamma}{N}$, the condition (2.10) can be simplified to

$$
\begin{equation*}
P_{i} \succ \rho\left(\frac{N}{2-\gamma}-1\right) A_{i}^{\top} A_{i}, i=1,2, \ldots, N . \tag{2.11}
\end{equation*}
$$

For special choices of $P_{i}$ :

- $P_{i}=\tau_{i} \mathbf{I}$ (standard proximal), condition (2.11) reduces to $\tau_{i}>\rho\left(\frac{N}{2-\gamma}-1\right)\left\|A_{i}\right\|^{2} ;$
- $P_{i}=\tau_{i} \mathbf{I}-\rho A_{i}^{\top} A_{i}$ (prox-linear), condition (2.11) reduces to $\tau_{i}>\frac{\rho N}{2-\gamma}\left\|A_{i}\right\|^{2}$.

Note that one can normalize $A_{i}$ so that $\left\|A_{i}\right\| \leq 1$ and further simplify the above sufficient conditions.

Proof For any $\mathbf{u}=(\mathbf{x} ; \lambda) \in \mathbb{R}^{n+m}$, we have

$$
\begin{equation*}
\|\mathbf{u}\|_{Q}^{2}:=\|\mathbf{x}\|_{G_{x}}^{2}+\frac{2-\gamma}{\rho \gamma^{2}}\|\lambda\|^{2}+\frac{2}{\gamma} \lambda^{\top} A \mathbf{x} . \tag{2.12}
\end{equation*}
$$

Using the following basic inequality:

$$
\begin{equation*}
\frac{2}{\gamma} \lambda^{\top} A \mathbf{x}=\sum_{i=1}^{N} \frac{2}{\gamma} \lambda^{\top} A_{i} \mathbf{x}_{i} \geq-\sum_{i=1}^{N}\left(\frac{\epsilon_{i}}{\rho \gamma^{2}}\|\lambda\|^{2}+\frac{\rho}{\epsilon_{i}}\left\|A_{i} \mathbf{x}_{i}\right\|^{2}\right), \tag{2.13}
\end{equation*}
$$

for any $\epsilon_{i}>0(i=1,2, \ldots, N)$, we have

$$
\begin{equation*}
\|\mathbf{u}\|_{Q}^{2} \geq \sum_{i=1}^{N}\left\|\mathbf{x}_{i}\right\|_{P_{i}+\rho A_{i}^{\top} A_{i}-\frac{\rho}{\epsilon_{i}} A_{i}^{\top} A_{i}}^{2}+\frac{2-\gamma-\sum_{i=1}^{N} \epsilon_{i}}{\rho \gamma^{2}}\|\lambda\|^{2} . \tag{2.14}
\end{equation*}
$$

The condition (2.10) guarantees that $P_{i}+\rho A_{i}^{\top} A_{i}-\frac{\rho}{\epsilon_{i}} A_{i}^{\top} A_{i} \succ 0$ and $2-\gamma-\sum_{i=1}^{N} \epsilon_{i}>0$, and thus $\|\mathbf{u}\|_{Q}>0$. Hence, $Q$ is positive definite. The rest follows immediately.

### 2.2 Rate of Convergence

Next, we shall establish the $o(1 / k)$ convergence rate of Jacobi-Proximal ADMM. We use the quantity $\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{M}^{2}$ as a measure of the convergence rate motivated by [21,24]. Here, we define the matrix $M$ by

$$
M:=\left(\begin{array}{cc}
M_{x} & \\
& \frac{1}{\gamma \rho} \mathbf{I}
\end{array}\right) \quad \text { and } \quad M_{x}:=G_{x}-\rho A^{\top} A .
$$

Theorem 2.2 If $Q \succ 0$ and $M_{x} \succeq 0$, then $\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{M}^{2}=o(1 / k)$. Hence, $\left\|\mathbf{x}^{k}-\mathbf{x}^{k+1}\right\|_{M_{x}}^{2}=$ $o(1 / k)$ and $\left\|\lambda^{k}-\lambda^{k+1}\right\|^{2}=o(1 / k)$.

We need the following monotonic property of the iterations:
Lemma 2.2 If $M_{x} \succeq 0$ and $0<\gamma<2$, then $\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{M}^{2} \leq\left\|\mathbf{u}^{k-1}-\mathbf{u}^{k}\right\|_{M}^{2}$.
Proof Let $\Delta \mathbf{x}_{i}^{k+1}=\mathbf{x}_{i}^{k}-\mathbf{x}_{i}^{k+1}, i=1, \ldots, N, \Delta \mathbf{x}^{k+1}=\mathbf{x}^{k}-\mathbf{x}^{k+1}$, and $\Delta \lambda^{k+1}=\lambda^{k}-\lambda^{k+1}$. By Lemma 1.1, the optimality conditions (2.4) at $k$-th and $(k+1)$-th iterations yield

$$
\begin{aligned}
& \left\langle A_{i} \Delta \mathbf{x}_{i}^{k+1}, \Delta \lambda^{k}-\rho A \Delta \mathbf{x}^{k+1}-\rho \sum_{j \neq i} A_{j}\left(\Delta \mathbf{x}_{j}^{k}-\Delta \mathbf{x}_{j}^{k+1}\right)\right\rangle \\
& \quad+\left(\Delta \mathbf{x}_{i}^{k+1}\right)^{\top} P_{i}\left(\Delta \mathbf{x}_{i}^{k}-\Delta \mathbf{x}_{i}^{k+1}\right) \geq 0
\end{aligned}
$$

Summing up over all $i$ and rearranging the terms, we have

$$
\begin{equation*}
\left\langle A \Delta \mathbf{x}^{k+1}, \Delta \lambda^{k}\right\rangle \geq\left\|\Delta \mathbf{x}^{k+1}\right\|_{G_{x}}^{2}-\left(\Delta \mathbf{x}^{k}\right)^{\top}\left(G_{x}-\rho A^{\top} A\right) \Delta \mathbf{x}^{k+1} \tag{2.15}
\end{equation*}
$$

Since $M_{x}:=G_{x}-\rho A^{\top} A \succeq 0$, we have $2\left(\Delta \mathbf{x}^{k}\right)^{\top}\left(G_{x}-\rho A^{\top} A\right) \Delta \mathbf{x}^{k+1} \leq\left\|\Delta \mathbf{x}^{k}\right\|_{M_{x}}^{2}+$ $\left\|\Delta \mathbf{x}^{k+1}\right\|_{M_{x}}^{2}$, and thus

$$
2\left\langle A \Delta \mathbf{x}^{k+1}, \Delta \lambda^{k}\right\rangle \geq\left\|\Delta \mathbf{x}^{k+1}\right\|_{2 G_{x}-M_{x}}^{2}-\left\|\Delta \mathbf{x}^{k}\right\|_{M_{x}}^{2}=\left\|\Delta \mathbf{x}^{k+1}\right\|_{G_{x}+\rho A^{\top} A}^{2}-\left\|\Delta \mathbf{x}^{k}\right\|_{M_{x}}^{2}
$$

Note that $\Delta \lambda^{k+1}=\Delta \lambda^{k}-\gamma \rho A \Delta \mathbf{x}^{k+1}$. It follows that

$$
\begin{aligned}
\frac{1}{\gamma \rho}\left\|\Delta \lambda^{k}\right\|^{2}-\frac{1}{\gamma \rho}\left\|\Delta \lambda^{k+1}\right\|^{2} & =2\left\langle A \Delta \mathbf{x}^{k+1}, \Delta \lambda^{k}\right\rangle-\gamma \rho\left\|A \Delta \mathbf{x}^{k+1}\right\|^{2} \\
& \geq\left\|\Delta \mathbf{x}^{k+1}\right\|_{G_{x}+(1-\gamma) \rho A^{\top} A}^{2}-\left\|\Delta \mathbf{x}^{k}\right\|_{M_{x}}^{2}
\end{aligned}
$$

i.e.,
$\left(\left\|\Delta \mathbf{x}^{k}\right\|_{M_{x}}^{2}+\frac{1}{\gamma \rho}\left\|\Delta \lambda^{k}\right\|^{2}\right)-\left(\left\|\Delta \mathbf{x}^{k+1}\right\|_{M_{x}}+\frac{1}{\gamma \rho}\left\|\Delta \lambda^{k+1}\right\|^{2}\right) \geq\left\|\Delta \mathbf{x}^{k+1}\right\|_{(2-\gamma) \rho A^{\top} A}^{2} \geq 0$, which completes the proof.

Proof of Theorem 2.2 By Theorem 2.1, there exists some $\eta>0$ such that

$$
\begin{equation*}
\left\|\mathbf{u}^{k}-\mathbf{u}^{*}\right\|_{G}^{2}-\left\|\mathbf{u}^{k+1}-\mathbf{u}^{*}\right\|_{G}^{2} \geq\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{Q}^{2} \geq \eta\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{M}^{2} \tag{2.16}
\end{equation*}
$$

Summing (2.16) over $k$ gives $\sum_{k=1}^{\infty}\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{M}^{2}<\infty$. On the other hand, Lemma 2.2 implies the monotone non-increasing of $\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{M}^{2}$. By Lemma 1.1, we have $\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{M}^{2}=o(1 / k)$.

### 2.3 Adaptive Parameter Tuning

The parameters satisfying the condition (2.10) may be rather conservative, because the inequality (2.13) for bounding $\|\mathbf{u}\|_{Q}^{2}$ is usually very loose. In practice, we can compute $\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{Q}^{2}$ exactly at very little extra cost. If $\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{Q}^{2}>0$, then Lemma 2.1 assures the decreasing of the solution error (in the $G$-norm) so that the current parameters are acceptable. On the other hand, if $\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{Q}^{2}<0$, then the matrix $Q$ is not positive definite,
meaning that the current parameters $P_{i}, i=1,2, \ldots, N$ may be too small. So we should make $\left\{P_{i}\right\}$ bigger until $\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{Q}^{2}>0$ holds. Therefore, we propose a practical strategy for adaptively adjusting the matrices $\left\{P_{i}\right\}$ based on the value of $\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{Q}^{2}$ :

```
Initialize with small \(P_{i}^{0} \succeq 0(i=1,2, \ldots, N)\) and a small \(\eta>0\);
for \(k=1,2, \ldots\) do
    if \(\left\|\mathbf{u}^{k-1}-\mathbf{u}^{k}\right\|_{Q}^{2}>\eta \cdot\left\|\mathbf{u}^{k-1}-\mathbf{u}^{k}\right\|^{2}\) then
        \(P_{i}^{k+1} \leftarrow P_{i}^{k}, \forall i ;\)
    else
        Increase \(P_{i}: P_{i}^{k+1} \leftarrow \alpha_{i} P_{i}^{k}+\beta_{i} Q_{i}\left(\alpha_{i}>1, \beta_{i} \geq 0, Q_{i} \succ 0\right), \forall i\);
        Restart: \(\mathbf{u}^{k} \leftarrow \mathbf{u}^{k-1}\);
```

The above strategy starts with relatively small proximal parameters $\left\{P_{i}\right\}$ and gradually increase them. By Theorem 2.1, we know that when the parameters $\left\{P_{i}\right\}$ are large enough for (2.10) to hold, the condition (2.8) will be satisfied (for sufficiently small $\eta$ ). Therefore, the adjustment of $\left\{P_{i}\right\}$ cannot occur infinite times. After a finite number of iterations, $\left\{P_{i}\right\}$ will remain constant and the contraction property (2.9) of the iterations will hold. Therefore, the convergence of such an adaptive parameter tuning scheme follows immediately from our previous analysis.

Theorem 2.3 When the matrices $P_{i}(i=1,2, \ldots, N)$ in Algorithm 4 are adaptively adjusted in the above scheme, the algorithm converges to a solution to the problem (1.1).

Empirical evidence shows that the parameters $\left\{P_{i}\right\}$ typically adjust themselves only during the first few iterations and then stay constant. Alternatively, one may also decrease the parameters after every few iterations or after they have not been updated for a certain number of iterations. But the total times of decrease should be bounded to guarantee convergence. By using this adaptive strategy, the resulting parameters $\left\{P_{i}\right\}$ are usually much smaller than those required by the condition (2.10), thereby leading to substantially faster convergence in practice.

## 3 Numerical Experiments

In this section, we present numerical results to compare the following parallel splitting algorithms:

- Prox-JADMM proposed Jacobi-Proximal ADMM (Algorithm 4);
- VSADMM Variable Splitting ADMM (Algorithm 1);
- Corr-JADMM Jacobi ADMM with correction steps [21]. At every iteration, it first generates a "predictor" $\tilde{\mathbf{u}}^{k+1}$ by an iteration of Jacobi ADMM (Algorithm 3) and then corrects $\tilde{\mathbf{u}}^{k+1}$ to generate the new iterate by:

$$
\begin{equation*}
\mathbf{u}^{k+1}=\mathbf{u}^{k}-\alpha_{k}\left(\mathbf{u}^{k}-\tilde{\mathbf{u}}^{k+1}\right), \tag{3.1}
\end{equation*}
$$

where $\alpha_{k}>0$ is a step size. In our experiments, we adopt the dynamically updated step size $\alpha_{k}$ according to [21], which is shown to converge significantly faster than using a constant step size, though updating the step size requires extra computation.

- YALLI one of the state-of-the-art solvers for the $\ell_{1}$-minimization problem.

In Sects. 3.1 and 3.2, all of the numerical experiments are run in MATLAB (R2011b) on a workstation with an Intel Core i5-3570 CPUs ( 3.40 GHz ) and 32 GB of RAM. Section 3.3 gives two very large instances that are solved by a C/MPI implementation on Amazon Elastic Compute Cloud (EC2).

### 3.1 Exchange Problem

Consider a network of $N$ agents that exchange $n$ commodities. Let $\mathbf{x}_{i} \in \mathbb{R}^{n}(i=1,2, \ldots, N)$ denote the amount of commodities that are exchanged among the $N$ agents. Each agent $i$ has a certain cost function $f_{i}: \mathbb{R}^{n} \rightarrow \mathbb{R}$. The exchange problem (see, e.g., [3] for a review) is given by

$$
\begin{equation*}
\min _{\left\{\mathbf{x}_{i}\right\}} \sum_{i=1}^{N} f_{i}\left(\mathbf{x}_{i}\right) \text { s.t. } \sum_{i=1}^{N} \mathbf{x}_{i}=0 \tag{3.2}
\end{equation*}
$$

which minimizes the total cost among $N$ agents subject to an equilibrium constraint on the commodities. This is a special case of (1.1) where $A_{i}=\mathbf{I}$ and $c=0$.

We consider quadratic cost functions $f_{i}\left(\mathbf{x}_{i}\right):=\frac{1}{2}\left\|C_{i} \mathbf{x}_{i}-d_{i}\right\|^{2}$, where $C_{i} \in \mathbb{R}^{p \times n}$ and $d_{i} \in \mathbb{R}^{p}$. All the compared algorithms solve the following type of subproblem:

$$
\begin{equation*}
\mathbf{x}_{i}^{k+1}=\arg \min _{\mathbf{x}_{i}} \frac{1}{2}\left\|C_{i} \mathbf{x}_{i}-d_{i}\right\|^{2}+\frac{\rho}{2}\left\|\mathbf{x}_{i}-b_{i}^{k}\right\|^{2}, \forall i=1,2, \ldots, N, \tag{3.3}
\end{equation*}
$$

except that Prox-JADMM also adds a proximal term $\frac{1}{2}\left\|\mathbf{x}_{i}-\mathbf{x}_{i}^{k}\right\|_{P_{i}}^{2}$. Here $b_{i}^{k} \in \mathbb{R}^{m}$ is a vector independent of $\mathbf{x}_{i}$ and takes different forms in different algorithms. For Prox-JADMM, we simply set $P_{i}=\tau_{i} \mathbf{I}\left(\tau_{i}>0\right)$. Clearly, each $\mathbf{x}_{i}$-subproblem is a quadratic program that can be computed efficiently using various methods.

In our experiment, we randomly generate $\mathbf{x}_{i}^{*}, i=1,2, \ldots, N-1$, following the standard Gaussian distribution, and let $\mathbf{x}_{N}^{*}=-\sum_{i=1}^{N-1} \mathbf{x}_{i}^{*}$. Matrices $C_{i}$ are random Gaussian matrices, and vectors $d_{i}$ are computed by $d_{i}=C_{i} \mathbf{x}_{i}^{*}$. Apparently, $\mathbf{x}^{*}$ is a solution (not necessarily unique) to (3.2), and the optimal objective value is 0 .

The penalty parameter $\rho$ is set to be $0.01,1$ and 0.01 for Prox-JADMM, VSADMM and Corr-JADMM, respectively. They are nearly optimal for each algorithm, picked out of a number of different values. Note that the parameter for VSADMM is quite different from the other two algorithms because it has different constraints due to the variable splitting. For Prox-JADMM, the proximal parameters are initialized by $\tau_{i}=0.1(N-1) \rho$ and adaptively updated by the strategy in Sect. 2.3; the parameter $\gamma$ is set to be 1 .

The size of the test problem is set to be $n=100, N=100, p=80$. Letting all the algorithms run 200 iterations, we plot their objective value $\sum_{i=1}^{N} f_{i}\left(\mathbf{x}_{i}\right)$ and residual $\left\|\sum_{i=1}^{N} \mathbf{x}_{i}\right\|_{2}$. Note that the per-iteration cost (both computation and communication) is roughly the same for all the compared algorithms. Figure 1 shows the comparison result, which is averaged over 100 random trials. We can see that Prox-JADMM is clearly the fastest one among the compared algorithms.

## $3.2 \ell_{1}$-Minimization

We consider the $\ell_{1}$-minimization problem for finding sparse solutions of an underdetermined linear system:

$$
\begin{equation*}
\min _{\mathbf{x}}\|\mathbf{x}\|_{1} \quad \text { s.t. } A \mathbf{x}=c, \tag{3.4}
\end{equation*}
$$



Fig. 1 Exchange problem $(n=100, N=100, p=80)$
where $\mathbf{x} \in \mathbb{R}^{n}, A \in \mathbb{R}^{m \times n}$ and $c \in \mathbb{R}^{m}(m<n)$. It is also known as the basis pursuit problem, which has been widely used in compressive sensing, signal and image processing, statistics, and machine learning. Suppose that the data is partitioned into $N$ blocks: $\mathbf{x}=\left[\mathbf{x}_{1}, \mathbf{x}_{2}, \ldots, \mathbf{x}_{N}\right]$ and $A=\left[A_{1}, A_{2}, \ldots, A_{N}\right]$. Then the problem (3.4) can be written in the form of (1.1) with $f_{i}\left(\mathbf{x}_{i}\right)=\left\|\mathbf{x}_{i}\right\|_{1}$.

In our experiment, a sparse solution $\mathbf{x}^{*}$ is randomly generated with $k(k \ll n)$ nonzeros drawn from the standard Gaussian distribution. Matrix $A$ is also randomly generated from the standard Gaussian distribution, and it is partitioned evenly into $N$ blocks. The vector $c$ is then computed by $c=A \mathbf{x}^{*}+\eta$, where $\eta \sim \mathcal{N}\left(0, \sigma^{2} \mathbf{I}\right)$ is Gaussian noise with standard deviation $\sigma$.

Prox-JADMM solves the $\mathbf{x}_{i}$-subproblems with $P_{i}=\tau_{i} \mathbf{I}-\rho A_{i}^{\top} A_{i}(i=1,2, \ldots, N)$ as follows:

$$
\begin{align*}
\mathbf{x}_{i}^{k+1} & =\arg \min _{\mathbf{x}_{i}}\left\|\mathbf{x}_{i}\right\|_{1}+\frac{\rho}{2}\left\|A_{i} \mathbf{x}_{i}+\sum_{j \neq i} A_{j} \mathbf{x}_{j}^{k}-c-\frac{\lambda^{k}}{\rho}\right\|^{2}+\frac{1}{2}\left\|\mathbf{x}_{i}-\mathbf{x}_{i}^{k}\right\|_{P_{i}}^{2} \\
& =\arg \min _{\mathbf{x}_{i}}\left\|\mathbf{x}_{i}\right\|_{1}+\left\langle\rho A_{i}^{\top}\left(A \mathbf{x}^{k}-c-\frac{\lambda^{k}}{\rho}\right), \mathbf{x}_{i}\right\rangle+\frac{\tau_{i}}{2}\left\|\mathbf{x}_{i}-\mathbf{x}_{i}^{k}\right\|^{2} . \tag{3.5}
\end{align*}
$$

Here, we choose the prox-linear $P_{i}$ 's to linearize the original subproblems, and thus (3.5) admits a simple closed-form solution by the shrinkage (or soft-thresholding) formula. The proximal parameters are initialized as $\tau_{i}=0.1 N \rho$ and are adaptively updated by the strategy discussed in Sect. 2.3.

Recall that VSADMM needs to solve the following $\mathbf{x}_{i}$-subproblems:

$$
\begin{equation*}
\mathbf{x}_{i}^{k+1}=\arg \min _{\mathbf{x}_{i}}\left\|\mathbf{x}_{i}\right\|_{1}+\frac{\rho}{2}\left\|A_{i} \mathbf{x}_{i}-\mathbf{z}_{i}^{k+1}-\frac{c}{N}-\frac{\lambda_{i}^{k}}{\rho}\right\|^{2} . \tag{3.6}
\end{equation*}
$$

Such subproblems are not easily computable, unless $\mathbf{x}_{i}$ is a scalar (i.e., $n_{i}=1$ ) or $A_{i}^{\top} A_{i}$ is a diagonal matrix. Instead, we solve the subproblems approximately using the prox-linear approach:
$\mathbf{x}_{i}^{k+1}=\arg \min _{\mathbf{x}_{i}}\left\|\mathbf{x}_{i}\right\|_{1}+\left\langle\rho A_{i}^{\top}\left(A_{i} \mathbf{x}_{i}^{k}-\mathbf{z}_{i}^{k+1}-\frac{c}{N}-\frac{\lambda_{i}^{k}}{\rho}\right), \mathbf{x}_{i}-\mathbf{x}_{i}^{k}\right\rangle+\frac{\tau_{i}}{2}\left\|\mathbf{x}_{i}-\mathbf{x}_{i}^{k}\right\|^{2}$,
which can be easily computed by the shrinkage operator. We set $\tau_{i}=1.01 \rho\left\|A_{i}\right\|^{2}$ in order to guarantee the convergence, as suggested in [36].

Corr-JADMM solves the following $\mathbf{x}_{i}$-subproblems in the "prediction" step:

$$
\begin{equation*}
\tilde{\mathbf{x}}_{i}^{k+1}=\arg \min _{\mathbf{x}_{i}}\left\|\mathbf{x}_{i}\right\|_{1}+\frac{\rho}{2}\left\|A_{i} \mathbf{x}_{i}+\sum_{j \neq i} A_{j} \mathbf{x}_{j}^{k}-c-\frac{\lambda^{k}}{\rho}\right\|^{2} \tag{3.7}
\end{equation*}
$$

Because the correction step in [21] is based on exact minimization of the subproblems, we do not apply the prox-linear approach to solve the subproblems approximately. Instead, we always partition $\mathbf{x}$ into scalar components (i.e., $N=n$ ) so that the subproblems (3.7) can still be computed exactly. The same penalty parameter $\rho=10 /\|c\|_{1}$ is used for the three algorithms. It is nearly optimal for each algorithm, selected out of a number of different values.

We also include the YALL1 package [37] in the experiment, which is one of the state-of-the-art solvers for $\ell_{1}$ minimization. Though YALL1 is not implemented in parallel, the matrix-vector multiplication by $A$ and $A^{\top}$ in the code can be parallelized (see [32]). Since all the compared algorithms have roughly the same amount of per-iteration cost (in terms of both computation and communication), we simply let all the algorithms run for a fixed number of iterations and plot their relative error $\frac{\left\|\mathbf{x}^{k}-\mathbf{x}^{*}\right\|_{2}}{\left\|\mathbf{x}^{*}\right\|_{2}}$.

Figure 2 shows the comparison result where $n=1000, m=300, k=60$ and the standard deviation of noise $\sigma$ is set to be 0 and $10^{-3}$, respectively. For Prox-JADMM and VSADMM, we set $N=100$; for Corr-JADMM, we set $N=1000$. The results are average of 100 random trials. We can see that Prox-JADMM and Corr-JADMM achieve very close performance and are the fastest ones among the compared algorithms. YALL1 also shows competitive performance. However, VSADMM is far slower than the others, probably due to inexact minimization of the subproblems and the conservative proximal parameters.



Fig. $2 \ell_{1}$-problem $(n=1000, m=300, k=60)$. a Noise-free $(\sigma=0)$, $\mathbf{b}$ noise added $\left(\sigma=10^{-3}\right)$

Table 1 Two large datasets

|  | $m$ | $n$ | RAM |
| :--- | :--- | :--- | :--- |
| dataset 1 | $1.0 \times 10^{5}$ | $2.0 \times 10^{5}$ | 150 GB |
| dataset 2 | $1.5 \times 10^{5}$ | $3.0 \times 10^{5}$ | 337 GB |

### 3.3 Distributed Large-Scale $\ell_{1}$-Minimization

We now use Algorithm 4 to solve two very large instances of the $\ell_{1}$-minimization problem (3.4) using a C code with MPI (for inter-process communication) and the GNU Scientific Library (GSL) (for BLAS operations). The experiments were carried out on Amazon EC2.

We generate two test instances as shown in Table 1. Specifically, a sparse solution $\mathbf{x}^{*}$ is randomly generated with $5 \%$ and $15 \%$ nonzeros drawn from the standard Gaussian distribution. Matrix $A$ is also randomly generated from the standard Gaussian distribution with $m$ rows and $n$ columns, and it is partitioned evenly into $N=80$ blocks. Vector $c$ is then computed by $c=A \mathbf{x}^{*}$. Note that $A$ is dense and has double precision. For Test 1 it requires over 150 GB of RAM and has 20 billion nonzero entries, and for Test 2 it requires over 337 GB of RAM. Those two tests are far too large to process on a single PC or workstation. We want to point out that we cannot find a dataset of similar or larger size in the public domain. We are willing to try a larger problem per reader's request. Note that this is the first time that an alternating direction type of algorithm has ever been tested on problems of such a large scale.

We solve the problem on a cluster of 10 machines, where each machine is a "memoryoptimized instance" with 68 GB RAM and 1 eight-core Intel Xeon E5-2665 CPU. They run Ubuntu 12.04 and are connected with 10 Gigabit ethernet network. Since each has 8 cores, we run the code with 80 processes so that each process runs on its own core. Such a setup is charged for under $\$ 17$ per hour.

We solve the large-scale $\ell_{1}$ minimization problems with a C code that matches our Matlab code. The C code has about 300 lines and it is available for download on authors' website. ${ }^{1}$

The breakdown of the wall-clock time is summarized in Tables 2 and 3. We can observe that Jacobi ADMM is very efficient in obtaining a relative low accuracy, which is usually sufficient for large-scale problems. We want to point out that the basic BLAS operations in our implantation can be further improved by using other libraries such as hardware-optimized BLAS libraries produced by ATLAS, Armadillo, etc. Those libraries might lead to several times of speedup. ${ }^{2}$ We use GSL due to its ease of use and easy adaptation to other problems.

## 4 A Sufficient Condition for Convergence of Jacobi ADMM

In this section, we provide a sufficient condition to guarantee the convergence of Jacobi ADMM (Algorithm 3), which does not use proximal terms or any correction step. The condition only depends on the matrices $A_{i}$, without imposing further assumptions on the objective functions $f_{i}$ or the penalty parameter $\rho$. For the Gauss-Seidel ADMM (Algorithm 2), a sufficient condition for convergence is provided in [5] for the special case $N=3$, assuming two of the three coefficient matrices are orthogonal. Our condition does not require

[^0]Table 2 Time results for large scale $\ell_{1}$ minimization examples (sparsity $5 \%$ )

|  | 150 GB test |  |  |  |  |  |  |  | 337 GB test |  |  |  |  |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | Itr | Time $(\mathrm{m})$ | Cost $(\$)$ | Itr | Time $(\mathrm{m})$ | Cost $(\$)$ |  |  |  |  |  |  |  |
| Data generation | - | 0.7 | 0.2 | - | 1.6 | 0.5 |  |  |  |  |  |  |  |
| CPU per iteration | - | 1.3 | - | - | 2.9 | - |  |  |  |  |  |  |  |
| Comm. per iteration | - | 0.07 | - | - | 0.15 | - |  |  |  |  |  |  |  |
| Reach $10^{-1}$ | 46 | 1.0 | 0.2 | 56 | 2.8 | 0.4 |  |  |  |  |  |  |  |
| Reach $10^{-2}$ | 176 | 3.9 | 0.6 | 203 | 10.0 | 1.6 |  |  |  |  |  |  |  |
| Reach $10^{-3}$ | 791 | 17.5 | 2.9 | 880 | 43.5 | 7.1 |  |  |  |  |  |  |  |

Table 3 Time results for large scale $\ell_{1}$ minimization examples (sparsity 15\%)

|  | 150 GB test |  | 337 GB test |  |  |  |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- |
|  | Itr | Time $(\mathrm{m})$ | Cost $(\$)$ | Itr | Time $(\mathrm{m})$ | Cost $(\$)$ |
| Reach $10^{-1}$ | 115 | 2.5 | 0.4 | 181 | 8.8 | 1.4 |
| Reach $10^{-2}$ | 614 | 13.4 | 2.1 | 745 | 36.2 | 5.9 |
| Reach $10^{-3}$ | 1000 | 21.8 | 3.6 | 1000 | 48.5 | 7.9 |
|  | $\left(\right.$ reach $\left.5.0 \times 10^{-3}\right)$ |  |  | $\left(\right.$ reach $\left.6.0 \times 10^{-3}\right)$ |  |  |

exact orthogonality. Instead, we mainly assume that the matrices $A_{i}, i=1,2, \ldots, N$ are mutually "near-orthogonal" and have full column-rank.

Theorem 4.1 If there exists $\delta \geq 0$ such that

$$
\begin{equation*}
\left\|A_{i}^{\top} A_{j}\right\| \leq \delta, \forall i \neq j, \text { and } \lambda_{\min }\left(A_{i}^{\top} A_{i}\right)>3(N-1) \delta, \forall i, \tag{4.1}
\end{equation*}
$$

where $\lambda_{\text {min }}\left(A_{i}^{\top} A_{i}\right)$ denotes the smallest eigenvalue of $A_{i}^{\top} A_{i}$, then the sequence $\left\{\mathbf{u}^{k}\right\}$ generated by Algorithm 3 converges to a solution $\mathbf{u}^{*}$ to the problem (1.1).

The proof technique is motivated by the contraction analysis of the sequence $\left\{\mathbf{u}^{k}\right\}$ under some $G$-norm (e.g., $[12,21,23]$ ). We first need the following lemma:

## Lemma 4.1 Let

$$
G_{0}:=\left(\begin{array}{cccc}
\rho A_{1}^{\top} A_{1} & & & \\
& \ddots & & \\
& & \rho A_{N}^{\top} A_{N} & \\
& & & \\
& \frac{1}{\rho} \mathbf{I}
\end{array}\right), Q_{0}:=\left(\begin{array}{ccccc}
\rho A_{1}^{\top} A_{1} & & & A_{1}^{\top} \\
& \ddots & & \vdots \\
& & \rho A_{N}^{\top} A_{N} & A_{N}^{\top} \\
& A_{1} & \ldots & A_{N} & \frac{1}{\rho} \mathbf{I}
\end{array}\right),
$$

where $\mathbf{I}$ is the identity matrix of size $m \times m$. For $k \geq 1$, we have

$$
\begin{equation*}
\left\|\mathbf{u}^{k}-\mathbf{u}^{*}\right\|_{G_{0}}^{2}-\left\|\mathbf{u}^{k+1}-\mathbf{u}^{*}\right\|_{G_{0}}^{2} \geq\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{Q_{0}}^{2} \tag{4.2}
\end{equation*}
$$

where

$$
\begin{equation*}
\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{Q_{0}}^{2}:=\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{G_{0}}^{2}+2\left(\lambda^{k}-\lambda^{k+1}\right)^{\top} A\left(\mathbf{x}^{k}-\mathbf{x}^{k+1}\right) . \tag{4.3}
\end{equation*}
$$

This lemma follows directly from Lemma 2.1 since it is a special case with $\gamma=1$ and $P_{i}=0, \forall i$. Now we are ready to prove the theorem.

Proof of Theorem 4.1 By the assumption $\left\|A_{i}^{\top} A_{j}\right\| \leq \delta, i \neq j$, we have

$$
\begin{equation*}
\left|\sum_{i \neq j}\left\langle A_{i} \mathbf{a}_{i}, A_{j} \mathbf{b}_{j}\right\rangle\right| \leq \sum_{i \neq j} \delta\left\|\mathbf{a}_{i}\right\|\left\|\mathbf{b}_{j}\right\| \leq \frac{\delta}{2}(N-1)\left(\|\mathbf{a}\|^{2}+\|\mathbf{b}\|^{2}\right), \quad \forall \mathbf{a}, \mathbf{b} \tag{4.4}
\end{equation*}
$$

To simplify the notation, we let $\mathbf{a}_{i}^{k}:=\mathbf{x}_{i}^{k}-\mathbf{x}_{i}^{*}, i=1,2, \ldots, N$. Note that $\lambda^{k}-\lambda^{k+1}=$ $\rho A \mathbf{a}^{k+1}, \mathbf{x}^{k}-\mathbf{x}^{k+1}=\mathbf{a}^{k}-\mathbf{a}^{k+1}$. Then, we can rewrite (4.3) as

$$
\begin{aligned}
\frac{1}{\rho}\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{Q_{0}}^{2} & =\sum_{i}\left\|A_{i}\left(\mathbf{a}_{i}^{k}-\mathbf{a}_{i}^{k+1}\right)\right\|^{2}+\left\|A \mathbf{a}^{k+1}\right\|^{2}+2\left\langle A \mathbf{a}^{k+1}, A\left(\mathbf{a}^{k}-\mathbf{a}^{k+1}\right)\right\rangle \\
& =\sum_{i}\left\|A_{i} \mathbf{a}_{i}^{k}\right\|^{2}+2 \sum_{i \neq j}\left\langle A_{i} \mathbf{a}_{i}^{k+1}, A_{j} \mathbf{a}_{j}^{k}\right\rangle-\sum_{i \neq j}\left\langle A_{i} \mathbf{a}_{i}^{k+1}, A_{j} \mathbf{a}_{j}^{k+1}\right\rangle \\
& \geq \sum_{i}\left\|A_{i} \mathbf{a}_{i}^{k}\right\|^{2}-(N-1) \delta\left(\left\|\mathbf{a}^{k+1}\right\|^{2}+\left\|\mathbf{a}^{k}\right\|^{2}\right)-(N-1) \delta\left\|\mathbf{a}^{k+1}\right\|^{2} \\
& =\sum_{i}\left\|A_{i} \mathbf{a}_{i}^{k}\right\|^{2}-(N-1) \delta\left\|\mathbf{a}^{k}\right\|^{2}-2(N-1) \delta\left\|\mathbf{a}^{k+1}\right\|^{2},
\end{aligned}
$$

where the inequality comes from (4.4). By Lemma 4.1, we have

$$
\begin{align*}
\left\|\mathbf{u}^{k}-\mathbf{u}^{*}\right\|_{G_{0}}^{2}-2(N-1) \delta \rho\left\|\mathbf{a}^{k}\right\|^{2} \geq & \left\|\mathbf{u}^{k+1}-\mathbf{u}^{*}\right\|_{G_{0}}^{2}-2(N-1) \delta \rho\left\|\mathbf{a}^{k+1}\right\|^{2} \\
& +\rho \sum_{i}\left\|A_{i} \mathbf{a}_{i}^{k}\right\|^{2}-3(N-1) \delta \rho\left\|\mathbf{a}^{k}\right\|^{2} . \tag{4.5}
\end{align*}
$$

We further simplify (4.5) as

$$
\begin{equation*}
b^{k}-b^{k+1} \geq d^{k} \tag{4.6}
\end{equation*}
$$

where the sequences $\left\{b^{k}\right\}$ and $\left\{d^{k}\right\}$ are defined by

$$
\begin{align*}
b^{k} & :=\left\|\mathbf{u}^{k}-\mathbf{u}^{*}\right\|_{G_{0}}^{2}-2(N-1) \delta \rho\left\|\mathbf{a}^{k}\right\|^{2},  \tag{4.7}\\
d^{k} & :=\rho \sum_{i}\left\|A_{i} \mathbf{a}_{i}^{k}\right\|^{2}-3(N-1) \delta \rho\left\|\mathbf{a}^{k}\right\|^{2} . \tag{4.8}
\end{align*}
$$

By the definition of $G_{0}$, we have

$$
\begin{equation*}
b^{k}=\rho \sum_{i}\left\|A_{i} \mathbf{a}_{i}^{k}\right\|^{2}-2(N-1) \delta \rho\left\|\mathbf{a}_{i}^{k}\right\|^{2}+\frac{1}{\rho}\left\|\lambda^{k}-\lambda^{*}\right\|^{2} . \tag{4.9}
\end{equation*}
$$

Since we assume $\lambda_{\min }\left(A_{i}^{\top} A_{i}\right)>3(N-1) \delta$, it follows that

$$
\begin{equation*}
\left\|A_{i} \mathbf{a}_{i}^{k}\right\|^{2} \geq 3(N-1) \delta\left\|\mathbf{a}_{i}^{k}\right\|^{2}, \forall i . \tag{4.10}
\end{equation*}
$$

Then it is easy to see that $b^{k} \geq 0$ and $d^{k} \geq 0$. By (4.6), the nonnegative sequence $\left\{b^{k}\right\}$ is monotonically non-increasing. Hence, $\left\{b^{k}\right\}$ converges to some $b^{*} \geq 0$. By (4.6), it also follows that $d^{k} \rightarrow 0$. Therefore, $\mathbf{a}^{k} \rightarrow 0$, i.e., $\mathbf{x}^{k} \rightarrow \mathbf{x}^{*}$.

Next we show $\lambda^{k} \rightarrow \lambda^{*}$. By taking limit of (4.9) and using $\mathbf{a}^{k} \rightarrow 0$, we have

$$
\begin{equation*}
b^{*}=\lim _{k \rightarrow \infty} b^{k}=\lim _{k \rightarrow \infty} \frac{1}{\rho}\left\|\lambda^{k}-\lambda^{*}\right\|^{2} . \tag{4.11}
\end{equation*}
$$

To show $\lambda^{k} \rightarrow \lambda^{*}$, it thus suffices to show $b^{*}=0$.
By (4.11), $\left\{\lambda^{k}\right\}$ is bounded and must have a convergent subsequence $\lambda^{k_{j}} \rightarrow \bar{\lambda}$. Recall the optimality conditions for the $\mathbf{x}_{i}$-subproblems (1.11):

$$
\begin{equation*}
A_{i}^{\top}\left(\lambda^{k}-\rho\left(A_{i} \mathbf{x}_{i}^{k+1}+\sum_{j \neq i} A_{j} \mathbf{x}_{j}^{k}-c\right)\right) \in \partial f_{i}\left(\mathbf{x}_{i}^{k+1}\right) . \tag{4.12}
\end{equation*}
$$

By Theorem 24.4 of [33], taking limit over the subsequence $\left\{k_{j}\right\}$ on both sides of (4.12) yields: $A_{i}^{\top} \bar{\lambda} \in \partial f_{i}\left(\mathbf{x}_{i}^{*}\right), \forall i$. Therefore, $\left(\mathbf{x}^{*}, \bar{\lambda}\right)$ satisfies the KKT conditions of the problem (1.1). Since ( $\mathbf{x}^{*}, \lambda^{*}$ ) is any KKT point, now we let $\lambda^{*}=\bar{\lambda}$. By (4.11) and $\left\|\lambda^{k_{j}}-\lambda^{*}\right\|^{2} \rightarrow 0$, we must have $b^{*}=0$, thereby completing the proof.

Similar to Theorem 2.2, we can find the convergence rate. That is,
Theorem 4.2 Let $W$ be blockly diagonal matrix $\operatorname{diag}\left(\rho \delta(N-1) I_{1}, \cdots, \rho \delta(N-1) I_{N}, \frac{1}{\rho} I\right)$. Under the assumptions in Theorem 4.1, $\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{W}^{2}=o(1 / k)$. That is, $\left\|\mathbf{x}^{k}-\mathbf{x}^{k+1}\right\|^{2}=$ $o(1 / k)$ and $\left\|\lambda^{k}-\lambda^{k+1}\right\|^{2}=o(1 / k)$.

Under the similar near-orthogonality assumption on the matrices $A_{i}, i=1,2, \ldots, N$, we have the following convergence result for Jacobi-Proximal ADMM:

Theorem 4.3 Suppose there exists $\delta \geq 0$ such that $\left\|A_{i}^{\top} A_{j}\right\| \leq \delta$ for all $i \neq j$, and the parameters in Algorithm 4 satisfy the following condition: for some $\alpha, \beta>0$ and $0<\gamma<2$,

$$
\left\{\begin{array}{l}
P_{i} \succ \rho\left(\frac{1}{\alpha}-1\right) A_{i}^{\top} A_{i}+\frac{\rho}{\beta} \delta(N-1) \mathbf{I}  \tag{4.13}\\
\lambda_{\min }\left(A_{i}^{\top} A_{i}\right)>\frac{2-\gamma+\beta}{2-\gamma-\alpha} \delta(N-1)
\end{array} \text { for } i=1, \ldots, N .\right.
$$

Then Algorithm 4 converges to a solution to the problem (1.1).
Proof Let

$$
H:=\left(\begin{array}{lll}
A_{1}^{\top} A_{1} & & \\
& \ddots & \\
& & A_{N}^{\top} A_{N}
\end{array}\right)
$$

If $\left\|A_{i}^{\top} A_{j}\right\| \leq \delta$ for all $i \neq j$, then it is easy to show the following: for any $\mathbf{x}$ and $\mathbf{y}$,

$$
\begin{aligned}
\|A \mathbf{x}\|^{2} & =\sum_{i=1}^{N}\left\|A_{i} \mathbf{x}_{i}\right\|^{2}+\sum_{i \neq j} \mathbf{x}_{i}^{\top} A_{i}^{\top} A_{j} \mathbf{x}_{j} \geq \sum_{i=1}^{N}\left\|A_{i} \mathbf{x}_{i}\right\|^{2}-\delta \sum_{i \neq j}\left\|\mathbf{x}_{i}\right\|\left\|\mathbf{x}_{j}\right\| \\
& \geq \sum_{i=1}^{N}\left\|A_{i} \mathbf{x}_{i}\right\|^{2}-\delta(N-1)\|\mathbf{x}\|^{2}=\|\mathbf{x}\|_{[H-\delta(N-1) \mathbf{I}]}^{2}
\end{aligned}
$$

and

$$
\begin{aligned}
2 \mathbf{x}^{\top} A^{\top} A \mathbf{y} & =2 \sum_{i=1}^{N} \mathbf{x}_{i}^{\top} A_{i}^{\top} A_{j} \mathbf{y}_{j}+2 \sum_{i \neq j} \mathbf{x}_{i}^{\top} A_{i}^{\top} A_{j} \mathbf{y}_{j} \geq 2 \sum_{i=1}^{N} \mathbf{x}_{i}^{\top} A_{i}^{\top} A_{j} \mathbf{y}_{j}-2 \delta \sum_{i \neq j}\left\|\mathbf{x}_{i}\right\|\left\|\mathbf{y}_{j}\right\| \\
& \geq-\sum_{i=1}^{N} \alpha\left\|A_{i} \mathbf{x}_{i}\right\|^{2}-\beta \delta(N-1)\|\mathbf{x}\|^{2}-\sum_{i=1}^{N} \frac{1}{\alpha}\left\|A_{i} \mathbf{y}_{i}\right\|^{2}-\frac{1}{\beta} \delta(N-1)\|\mathbf{y}\|^{2} \\
& =-\|\mathbf{x}\|_{[\alpha H+\beta \delta(N-1) \mathbf{I}]}^{2}-\|\mathbf{y}\|_{\left[\frac{1}{\alpha} H+\frac{1}{\beta} \delta(N-1) \mathbf{I}\right]}^{2}, \quad \forall \alpha, \beta>0,
\end{aligned}
$$

Using the above inequalities, we have

$$
\begin{aligned}
& \frac{2}{\gamma}\left(\lambda^{k}-\lambda^{k+1}\right)^{\top} A\left(\mathbf{x}^{k}-\mathbf{x}^{k+1}\right)=2 \rho\left(\mathbf{x}^{k+1}-\mathbf{x}^{*}\right) A^{\top} A\left(\mathbf{x}^{k}-\mathbf{x}^{k+1}\right) \\
\geq & -\rho\left\|\mathbf{x}^{k+1}-\mathbf{x}^{*}\right\|_{[\alpha H+\beta \delta(N-1) \mathbf{I}]}^{2}-\rho\left\|\mathbf{x}^{k}-\mathbf{x}^{k+1}\right\|_{\left[\frac{1}{\alpha} H+\frac{1}{\beta} \delta(N-1) \mathbf{I}\right]}^{2},
\end{aligned}
$$

and $\left\|\lambda^{k}-\lambda^{k+1}\right\|^{2}=\gamma^{2} \rho^{2}\left\|A\left(\mathbf{x}^{k+1}-\mathbf{x}^{*}\right)\right\|^{2} \geq \gamma^{2} \rho^{2}\left\|\mathbf{x}^{k+1}-\mathbf{x}^{*}\right\|_{[H-\delta(N-1) \mathbf{I}]}^{2}$. Therefore,

$$
\begin{aligned}
\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{Q}^{2} \geq & \left\|\mathbf{x}^{k}-\mathbf{x}^{k+1}\right\|_{G_{x}}^{2}+(2-\gamma) \rho\left\|\mathbf{x}^{k+1}-\mathbf{x}^{*}\right\|_{[H-\delta(N-1) \mathbf{I}]}^{2} \\
& -\rho\left\|\mathbf{x}^{k+1}-\mathbf{x}^{*}\right\|_{[\alpha H+\beta \delta(N-1) \mathbf{I}]}^{2}-\rho\left\|\mathbf{x}^{k}-\mathbf{x}^{k+1}\right\|_{\left[\frac{1}{\alpha} H+\frac{1}{\beta} \delta(N-1) \mathbf{I}\right]}^{2} .
\end{aligned}
$$

As long as the following holds:

$$
\left\{\begin{array}{l}
G_{x} \succ \frac{\rho}{\alpha} H+\frac{\rho}{\beta} \delta(N-1) \mathbf{I}, \\
(2-\gamma) \rho[H-\delta(N-1) \mathbf{I}] \succ \rho[\alpha H+\beta \delta(N-1) \mathbf{I}],
\end{array}\right.
$$

which is equivalent to the condition (4.13), there exists some $\eta>0$ such that (2.8) and (2.9) hold. Then the convergence of Algorithm 4 follows immediately from the standard analysis of contraction methods [19].

Remark 4.1 The conditions in Theorem 4.1 and Theorem 4.3 guarantee that (2.8) holds, i.e., there exists some $\eta>0$ such that $\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|_{Q}^{2} \geq \eta \cdot\left\|\mathbf{u}^{k}-\mathbf{u}^{k+1}\right\|^{2}$. But the matrix $Q$ (or $\left.Q_{0}\right)$ is not necessarily positive semi-definite. Also, the term $\frac{2-\gamma+\beta}{2-\gamma-\alpha} \delta(N-1)$ in (4.13) may be negative for some $\alpha$. Then, $\lambda_{\min }\left(A_{i}^{\top} A_{i}\right)$ is allowed to be 0 , in other words, $A_{i}$ may not be of full column rank. As long as the conditions in (4.13) are satisfied, Algorithm 4 will converge.

## 5 Connections between the Prox-JADMM Algorithm and the Proximal Point Algorithm

The Jacobi-Proximal ADMM is equivalent to the proximal point algorithm (PPA) applied to the following optimal condition of (1.1):

$$
\left[\begin{array}{c}
0  \tag{5.1}\\
0 \\
\vdots \\
0 \\
0
\end{array}\right] \in\left[\begin{array}{c}
\partial f_{1}\left(\mathbf{x}_{1}\right) \\
\partial f_{2}\left(\mathbf{x}_{2}\right) \\
\vdots \\
\partial f_{N}\left(\mathbf{x}_{N}\right) \\
0
\end{array}\right]+\underbrace{\left[\begin{array}{ccccc}
0 & 0 & \cdots & 0 & -A_{1}^{T} \\
0 & 0 & \cdots & 0 & -A_{2}^{T} \\
\vdots & \vdots & \vdots & \vdots & \vdots \\
0 & 0 & \cdots & 0 & -A_{N}^{T} \\
A_{1} & A_{2} & \cdots & A_{N} & 0
\end{array}\right]}_{B}\left[\begin{array}{c}
\mathbf{x}_{1} \\
\mathbf{x}_{2} \\
\vdots \\
\mathbf{x}_{N} \\
\lambda
\end{array}\right]-\left[\begin{array}{c}
0 \\
0 \\
\vdots \\
0 \\
c
\end{array}\right] .
$$

Note that $\mathbf{x}$ is an optimal solution if and only if $\mathbf{x}_{1}, \ldots, \mathbf{x}_{m}$, and $\lambda$ satisfy (5.1). Define $\mathbf{z}:=$ $\left(\mathbf{x}_{1}, \ldots, \mathbf{x}_{N}, \lambda\right), F(\mathbf{z}):=\sum_{i=1}^{N} f_{i}\left(\mathbf{x}_{i}\right)$, and let $S \mathbf{z}=B \mathbf{z}-c$, then finding $\mathbf{z}$ satisfying (5.1) is equivalent to

$$
\begin{equation*}
\text { findingzsuch that } 0 \in \partial F(\mathbf{z})+S \mathbf{z} \text {. } \tag{5.2}
\end{equation*}
$$

Introducing a symmetric positive definite matrix $W$, and applying PPA to (5.2) gives

$$
\begin{equation*}
\mathbf{z}^{k+1}=J_{W^{-1}(\partial F+S)}\left(\mathbf{z}^{k}\right) \tag{5.3}
\end{equation*}
$$

where $J_{A}:=(I+A)^{-1}$ is the resolvent of the operator $A$. Next, we show that (5.3) is equivalent to Prox-JADMM with a proper choice of $W$. Let

$$
W=\left[\begin{array}{cc}
Q-\rho A^{T} A & 0  \tag{5.4}\\
0 & \frac{1}{\rho} I
\end{array}\right],
$$

where $Q=\operatorname{diag}\left(Q_{1}, \ldots, Q_{N}\right), Q_{i}$ are symmetric matrices and $A=\left(A_{1}, \ldots, A_{N}\right)$. We require $Q \succ \rho A^{T} A$ so that $W$ is symmetric positive definite. The reason for choosing such a $W$ is to decouple $\mathbf{x}_{1}, \ldots, \mathbf{x}_{N}$ and $\lambda$. Based on the definition of the resolvent operator, we have

$$
\begin{equation*}
W \mathbf{z}^{k+1}=W \mathbf{z}^{k}-\tilde{\nabla} F\left(\mathbf{z}^{k+1}\right)-S \mathbf{z}^{k+1}, \tag{5.5}
\end{equation*}
$$

where $\tilde{\nabla} F\left(\mathbf{z}^{k+1}\right) \in \partial F\left(\mathbf{z}^{k+1}\right)$. Substituting (5.4) to (5.5) gives the following update

$$
\begin{align*}
& \tilde{\nabla} f_{i}\left(\mathbf{x}_{i}^{k+1}\right)+Q_{i}\left(\mathbf{x}_{i}^{k+1}-\mathbf{x}^{k}\right)+\rho A_{i}^{T}\left(A \mathbf{x}^{k}-\frac{\lambda^{k}}{\rho}-c\right)=0,  \tag{5.6}\\
& \lambda^{k+1}=\lambda^{k}-\rho\left(A \mathbf{x}^{k+1}-c\right) \tag{5.7}
\end{align*}
$$

where (5.6) corresponds to the optimal condition of the following problem

$$
\begin{equation*}
\mathbf{x}_{i}^{k+1}=\arg \min _{\mathbf{x}_{i}} f_{i}\left(\mathbf{x}_{i}\right)+\left\langle\rho A_{i}^{T}\left(A \mathbf{x}^{k}-\frac{\lambda^{k}}{\rho}-c\right), \mathbf{x}_{i}\right\rangle+\frac{1}{2}\left\|\mathbf{x}_{i}-\mathbf{x}_{i}^{k}\right\|_{Q_{i}}^{2} \tag{5.8}
\end{equation*}
$$

To show the equivalence, we distinguish two cases:

1. Let $Q_{i}=\tau_{i} I$, then (5.8) and (5.7) correspond to Prox-JADMM with $P_{i}=\tau I-\rho A_{i}^{T} A_{i}$.
2. Let $Q_{i}=\tau_{i} I+\rho A_{i}^{T} A_{i}$, then (5.8) is equivalent to

$$
\begin{aligned}
\mathbf{x}_{i}^{k+1} & =\arg \min _{\mathbf{x}_{i}} f_{i}\left(\mathbf{x}_{i}\right)+\left\langle\rho A_{i}^{T}\left(A \mathbf{x}^{k}-\frac{\lambda^{k}}{\rho}-c\right), \mathbf{x}_{i}\right\rangle+\frac{1}{2}\left\|\mathbf{x}_{i}-\mathbf{x}_{i}^{k}\right\|_{\tau_{i} I+\rho A_{i}^{T} A_{i}}^{2} \\
& =\arg \min _{\mathbf{x}_{i}} f_{i}\left(\mathbf{x}_{i}\right)+\frac{\rho}{2}\left\|A_{i} \mathbf{x}_{i}+\sum_{j \neq i} A_{i} \mathbf{x}_{i}^{k}-\frac{\lambda^{k}}{\rho}-c\right\|_{2}^{2}+\frac{\tau_{i}}{2}\left\|\mathbf{x}_{i}-\mathbf{x}_{i}^{k}\right\|^{2},
\end{aligned}
$$

which is the primal update of Prox-JADMM with $P_{i}=\tau_{i} I$.

## 6 Conclusion

Due to the dramatically increasing demand for dealing with big data, parallel and distributed computational methods are highly desirable. ADMM, as a versatile algorithmic tool, has proven to be very effective at solving many large-scale problems and well suited for
distributed computing. Yet, its parallelization still needs further investigation and improvement. This paper proposes a simple parallel and distributed ADMM for solving problems with separable structures. The algorithm framework introduces more flexibility for solving the subproblems due to the use of proximal terms $\left\|\mathbf{x}_{i}-\mathbf{x}_{i}^{k}\right\|_{P_{i}}^{2}$ with wisely chosen $P_{i}$. Its theoretical properties such as global convergence and an $o(1 / k)$ rate are established. Our numerical results demonstrate the efficiency of the proposed method in comparison with several existing algorithms. The code is available online for further studies.

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## Appendix: On $o(1 / k)$ Convergence Rate of ADMM

The convergence of the standard two-block ADMM has been long established in the literature [14, 16]. Its convergence rate has been actively studied; see [ $9,10,12,17,23-25,28]$ and the references therein. In the following, we briefly review the convergence analysis for ADMM ( $N=2$ ) and then improve the $O(1 / k)$ convergence rate established in [24] slightly to $o(1 / k)$ by using the same technique as in Sect. 2.2.

As suggested in [24], the quantity $\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2}$ can be used to measure the optimality of the iterations of ADMM, where

$$
\mathbf{w}:=\binom{\mathbf{x}_{2}}{\lambda}, H:=\left(\begin{array}{ccc}
\rho A_{2}^{\top} A_{2} & \\
& & \frac{1}{\rho} \mathbf{I}
\end{array}\right)
$$

and $\mathbf{I}$ is the identity matrix of size $m \times m$. Note that $\mathbf{x}_{1}$ is not part of $\mathbf{w}$ because $\mathbf{x}_{1}$ can be regarded as an intermediate variable in the iterations of ADMM, whereas $\left(\mathbf{x}_{2}, \lambda\right)$ are the essential variables [3]. In fact, if $\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2}=0$ then $\mathbf{w}^{k+1}$ is optimal. The reasons are as follows. Recall the subproblems of ADMM:

$$
\begin{align*}
& \mathbf{x}_{1}^{k+1}=\arg \min _{\mathbf{x}_{1}} f_{1}\left(\mathbf{x}_{1}\right)+\frac{\rho}{2}\left\|A_{1} \mathbf{x}_{1}+A_{2} \mathbf{x}_{2}^{k}-\lambda^{k} / \rho\right\|^{2},  \tag{6.1}\\
& \mathbf{x}_{2}^{k+1}=\arg \min _{\mathbf{x}_{2}} f_{2}\left(\mathbf{x}_{2}\right)+\frac{\rho}{2}\left\|A_{1} \mathbf{x}_{1}^{k+1}+A_{2} \mathbf{x}_{2}-\lambda^{k} / \rho\right\|^{2} . \tag{6.2}
\end{align*}
$$

By the formula for $\lambda^{k+1}$, their optimality conditions can be written as:

$$
\begin{align*}
A_{1}^{\top} \lambda^{k+1}-\rho A_{1}^{\top} A_{2}\left(\mathbf{x}_{2}^{k}-\mathbf{x}_{2}^{k+1}\right) & \in \partial f\left(\mathbf{x}_{1}^{k+1}\right)  \tag{6.3}\\
A_{2}^{\top} \lambda^{k+1} & \in \partial f_{2}\left(\mathbf{x}_{2}^{k+1}\right) \tag{6.4}
\end{align*}
$$

In comparison with the KKT conditions (1.14a) and (1.14b), we can see that $\mathbf{u}^{k+1}=$ $\left(\mathbf{x}_{1}^{k+1}, \mathbf{x}_{2}^{k+1}, \lambda^{k+1}\right)$ is a solution of (1.1) if and only if the following holds:

$$
\begin{array}{ll}
\mathbf{r}_{p}^{k+1}:=A_{1} \mathbf{x}_{1}^{k+1}+A_{2} \mathbf{x}_{2}^{k+1}-c=0 & \text { (primal feasibility) } \\
\mathbf{r}_{d}^{k+1}:=\rho A_{1}^{\top} A_{2}\left(\mathbf{x}_{2}^{k}-\mathbf{x}_{2}^{k+1}\right)=0 & \text { (dual feasibility) } \tag{6.6}
\end{array}
$$

By the update formula for $\lambda^{k+1}$, we can write $\mathbf{r}_{p}$ equivalently as

$$
\begin{equation*}
\mathbf{r}_{p}^{k+1}=\frac{1}{\rho}\left(\lambda^{k}-\lambda^{k+1}\right) . \tag{6.7}
\end{equation*}
$$

Clearly, if $\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2}=0$ then the optimality conditions (6.5) and (6.6) are satisfied, so $\mathbf{w}^{k+1}$ is a solution. On the other hand, if $\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2}$ is large, then $\mathbf{w}^{k+1}$ is likely to be far away from being a solution. Therefore, the quantity $\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2}$ can be viewed as a measure of the distance between the iteration $\mathbf{w}^{k+1}$ and the solution set. Furthermore, based on the variational inequality (1.15) and the variational characterization of the iterations of ADMM, it is reasonable to use the quadratic term $\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2}$ rather than $\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}$ to measure the convergence rate of ADMM (see [24] for more details).

The work [24] proves that $\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2}$ converges to zero at a rate of $O(1 / k)$. The key steps of the proof are to establish the following properties:

- the sequence $\left\{\mathbf{w}^{k}\right\}$ is contractive:

$$
\begin{equation*}
\left\|\mathbf{w}^{k}-\mathbf{w}^{*}\right\|_{H}^{2}-\left\|\mathbf{w}^{k+1}-\mathbf{w}^{*}\right\|_{H}^{2} \geq\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2}, \tag{6.8}
\end{equation*}
$$

- the sequence $\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2}$ is monotonically non-increasing:

$$
\begin{equation*}
\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2} \leq\left\|\mathbf{w}^{k-1}-\mathbf{w}^{k}\right\|_{H}^{2} . \tag{6.9}
\end{equation*}
$$

The contraction property (6.8) has been long established and its proof dates back to [14,16]. Inspired by [24], we provide a shorter proof for (6.9) than the one in [24].
$\operatorname{Proof}\left(\right.$ Proof of (6.9)) Let $\Delta \mathbf{x}_{i}^{k+1}=\mathbf{x}_{i}^{k}-\mathbf{x}_{i}^{k+1}$ and $\Delta \lambda^{k+1}=\lambda^{k}-\lambda^{k+1}$. By Lemma 1.1, i.e., (1.17), the optimality condition 6.3 at the $k$-th and $(k+1)$-th iterations yields: $\left\langle\Delta \mathbf{x}_{1}^{k+1}, \quad A_{1}^{\top} \Delta \lambda^{k+1}-\rho A_{1}^{\top} A_{2}\left(\Delta \mathbf{x}_{2}^{k}-\Delta \mathbf{x}_{2}^{k+1}\right)\right\rangle \geq 0$. Similarly for (6.4), we obtain $\left\langle\Delta \mathbf{x}_{2}^{k+1}, A_{2}^{\top} \Delta \lambda^{k+1}\right\rangle \geq 0$. Adding the above two inequalities together, we have

$$
\begin{equation*}
\left(A_{1} \Delta \mathbf{x}_{1}^{k+1}+A_{2} \Delta \mathbf{x}_{2}^{k+1}\right)^{\top} \Delta \lambda^{k+1}-\rho\left(A_{1} \Delta \mathbf{x}_{1}^{k+1}\right)^{\top} A_{2}\left(\Delta \mathbf{x}_{2}^{k}-\Delta \mathbf{x}_{2}^{k+1}\right) \geq 0 \tag{6.10}
\end{equation*}
$$

Using the equality according to (6.7):

$$
\begin{equation*}
A_{1} \Delta \mathbf{x}_{1}^{k+1}+A_{2} \Delta \mathbf{x}_{2}^{k+1}=\frac{1}{\rho}\left(\Delta \lambda^{k}-\Delta \lambda^{k+1}\right) \tag{6.11}
\end{equation*}
$$

(6.10) becomes $\frac{1}{\rho}\left(\Delta \lambda^{k}-\Delta \lambda^{k+1}\right)^{\top} \Delta \lambda^{k+1}-\left(\Delta \lambda^{k}-\Delta \lambda^{k+1}-\rho A_{2} \Delta \mathbf{x}_{2}^{k+1}\right)^{\top} A_{2}\left(\Delta \mathbf{x}_{2}^{k}\right.$ $\left.-\Delta \mathbf{x}_{2}^{k+1}\right) \geq 0$. After rearranging the terms, we get

$$
\begin{align*}
& \left(\sqrt{\rho} A_{2} \Delta \mathbf{x}_{2}^{k}+\frac{1}{\sqrt{\rho}} \Delta \lambda^{k}\right)^{\top}\left(\sqrt{\rho} A_{2} \Delta \mathbf{x}_{2}^{k+1}+\frac{1}{\sqrt{\rho}} \Delta \lambda^{k+1}\right)-\left(A_{2} \Delta \mathbf{x}_{2}^{k}\right)^{\top} \Delta \lambda^{k} \\
& -\left(A_{2} \Delta \mathbf{x}_{2}^{k+1}\right)^{\top} \Delta \lambda^{k+1} \\
& \geq \frac{1}{\rho}\left\|\Delta \lambda^{k+1}\right\|^{2}+\rho\left\|A_{2} \Delta \mathbf{x}_{2}^{k+1}\right\|^{2}=\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2} \tag{6.12}
\end{align*}
$$

By the Cauchy-Schwarz inequality, we have $\left(a_{1}+b_{1}\right)^{\top}\left(a_{2}+b_{2}\right) \leq\left(\left\|a_{1}+b_{1}\right\|^{2}+\| a_{2}+\right.$ $\left.b_{2} \|^{2}\right) / 2$, or equivalently, $\left(a_{1}+b_{1}\right)^{\top}\left(a_{2}+b_{2}\right)-a_{1}^{\top} b_{1}-a_{2}^{\top} b_{2} \leq\left(\left\|a_{1}\right\|^{2}+\left\|b_{1}\right\|^{2}+\left\|a_{2}\right\|^{2}+\right.$
$\left.\left\|b_{2}\right\|^{2}\right) / 2$. Applying this inequality to the left-hand side of (6.12), we have

$$
\begin{aligned}
\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2} & \leq\left(\rho\left\|A_{2} \Delta \mathbf{x}_{2}^{k}\right\|^{2}+\frac{1}{\rho}\left\|\Delta \lambda^{k}\right\|^{2}+\rho\left\|A_{2} \Delta \mathbf{x}_{2}^{k+1}\right\|^{2}+\frac{1}{\rho}\left\|\Delta \lambda^{k+1}\right\|^{2}\right) / 2 \\
& =\left(\left\|\mathbf{w}^{k-1}-\mathbf{w}^{k}\right\|_{H}^{2}+\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2}\right) / 2
\end{aligned}
$$

and thus (6.9) follows immediately.
We are now ready to improve the convergence rate from $O(1 / k)$ to $o(1 / k)$.
Theorem 6.1 The sequence $\left\{\mathbf{w}^{k}\right\}$ generated by Algorithm 2 (for $N=2$ ) converges to a solution $\mathbf{w}^{*}$ of problem (1.1) in the $H$-norm, i.e., $\left\|\mathbf{w}^{k}-\mathbf{w}^{*}\right\|_{H}^{2} \rightarrow 0$, and $\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2}=$ $o(1 / k)$. Therefore,

$$
\begin{equation*}
\left\|A_{1} \mathbf{x}_{1}^{k}-A_{1} \mathbf{x}_{1}^{k+1}\right\|^{2}+\left\|A_{2} \mathbf{x}_{2}^{k}-A_{2} \mathbf{x}_{2}^{k+1}\right\|^{2}+\left\|\lambda^{k}-\lambda^{k+1}\right\|^{2}=o(1 / k) \tag{6.13}
\end{equation*}
$$

Proof Using the contractive property of the sequence $\left\{\mathbf{w}^{k}\right\}$ (6.8) along with the optimality conditions, the convergence of $\left\|\mathbf{w}^{k}-\mathbf{w}^{*}\right\|_{H}^{2} \rightarrow 0$ follows from the standard analysis for contraction methods [19].

By (6.8), we have

$$
\begin{equation*}
\sum_{k=1}^{n}\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2} \leq\left\|\mathbf{w}^{1}-\mathbf{w}^{*}\right\|_{H}^{2}-\left\|\mathbf{w}^{n+1}-\mathbf{w}^{*}\right\|_{H}^{2}, \forall n . \tag{6.14}
\end{equation*}
$$

Therefore, $\sum_{k=1}^{\infty}\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2}<\infty$. By (6.9), $\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2}$ is monotonically nonincreasing and nonnegative. So Lemma 1.1 indicates that $\left\|\mathbf{w}^{k}-\mathbf{w}^{k+1}\right\|_{H}^{2}=o(1 / k)$, which further implies that $\left\|A_{2} \mathbf{x}_{2}^{k}-A_{2} \mathbf{x}_{2}^{k+1}\right\|^{2}=o(1 / k)$ and $\left\|\lambda^{k}-\lambda^{k+1}\right\|^{2}=o(1 / k)$. By (6.11), we also have $\left\|A_{1} \mathbf{x}_{1}^{k}-A_{1} \mathbf{x}_{1}^{k+1}\right\|^{2}=o(1 / k)$. Thus (6.13) follows immediately.

Remark 6.1 The proof technique based on Lemma 1.1 can be applied to improve some other existing convergence rates of $O(1 / k)$ (e.g., $[8,21])$ to $o(1 / k)$ as well.

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[^0]:    ${ }^{1} \mathrm{https}: / /$ github.com/ZhiminPeng/Jacobi-ADMM.
    $2 \mathrm{http}: / /$ nghiaho.com/? $\mathrm{p}=1726$.

