# On the Expected Convergence of Randomly Permuted ADMM 

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#### Abstract

The alternating direction method of multipliers (ADMM) is now widely used in many fields, and its convergence was proved when two blocks of variables are alternately updated. It is computationally beneficial to extend the ADMM directly to the case of a multi-block (multiple variable blocks) convex minimization problem. Unfortunately, such an extension fails to converge even when solving a simple square system of linear equations. In this paper, however, we prove that, if in each step one randomly and independently permutes the updating order of any given number of blocks followed by the regular multiplier update, the method will converge in expectation for solving the square system of linear equations. Our analysis of random permutation will also be of independent interest.


## 1 Introduction

Consider a convex minimization problem with a separable objective function and linear constraints:

$$
\begin{array}{cl}
\min & f_{1}\left(x_{1}\right)+\cdots+f_{n}\left(x_{n}\right) \\
\mathrm{s.t.} & A_{1} x_{1}+\cdots+A_{n} x_{n}=b  \tag{1}\\
& x_{i} \in \mathcal{X}_{i}, i=1, \ldots, n
\end{array}
$$

where $A_{i} \in \mathbb{R}^{N \times d_{i}}, b \in \mathbb{R}^{N \times 1}, \mathcal{X}_{i} \subseteq \mathbb{R}^{d_{i}}$ is a closed convex set, and $f_{i}: \mathbb{R}^{d_{i}} \rightarrow \mathbb{R}$ is a closed convex function, $i=1, \ldots, n$. Many machine learning and engineering problems can be cast into linearly-constrained optimization problems with two blocks (see [3] for many examples) or more than two blocks (e.g. linear programming, robust principal component analysis, composite regularizers for structured sparsity; see [5, 24] for more examples).
ADMM (Alternating Direction Method of Multipliers) was first proposed in [10] (see also [4]|8]) to solve problem (1) when there are only two blocks (i.e. $n=2$ ). In this 2 -block case, the augmented

$$
\begin{align*}
& \text { Lagrangian function of } \mathbb{1}) \text { is } \\
& \qquad \mathcal{L}\left(x_{1}, x_{2} ; \mu\right)=f_{1}\left(x_{1}\right)+f_{2}\left(x_{2}\right)-\mu^{T}\left(A_{1} x_{1}+A_{2} x_{2}-b\right)+\frac{\beta}{2}\left\|A_{1} x_{1}+A_{2} x_{2}-b\right\|^{2}, \tag{2}
\end{align*}
$$

where $\mu$ is the Lagrangian multiplier and $\beta>0$ is the penalty parameter. Each iteration of ADMM consists of a cyclic update (i.e. Gauss-Seidal type update) of primal variables $x_{1}, x_{2}$ and a dual ascent type update of $\mu$ :

$$
\left\{\begin{array}{l}
x_{1}^{k+1}=\arg \min _{x_{1} \in \mathcal{X}_{1}} \mathcal{L}\left(x_{1}, x_{2}^{k} ; \mu^{k}\right)  \tag{3}\\
x_{2}^{k+1}=\arg \min _{x_{2} \in \mathcal{X}_{2}} \mathcal{L}\left(x_{1}^{k+1}, x_{2} ; \mu^{k}\right) \\
\mu^{k+1}=\mu^{k}-\beta\left(A_{1} x_{1}^{k+1}+A_{2} x_{2}^{k+1}-b\right)
\end{array}\right.
$$

[^0]Due to the separable structure of the objective function, each subproblem only involves one $f_{i}, i \in\{1,2\}$, thus may be easier to solve. This feature enables the wide application of ADMM in signal processing and statistical learning where the objective function of the problem can usually be decomposed as the sum of the loss function and the regularizer; see [3] for a review. The convergence of 2-block ADMM has been well studied; see [7,9] for some recent reviews.

It is natural and computationally beneficial to extend the original 2-block ADMM directly to solve the general $n$-block problem (1):

$$
\left\{\begin{align*}
x_{1}^{k+1}= & \arg \min _{x_{1} \in \mathcal{X}_{1}} \mathcal{L}\left(x_{1}, x_{2}^{k}, \ldots, x_{n}^{k} ; \mu^{k}\right)  \tag{4}\\
& \vdots \\
x_{n}^{k+1}= & \arg \min _{x_{n} \in \mathcal{X}_{n}} \mathcal{L}\left(x_{1}^{k+1}, \ldots, x_{n-1}^{k+1}, x_{n} ; \mu^{k}\right) \\
\mu^{k+1}= & \mu^{k}-\beta\left(A_{1} x_{1}^{k+1}+\cdots+A_{n} x_{n}^{k+1}-b\right)
\end{align*}\right.
$$

where the augmented Lagrangian function

$$
\begin{equation*}
\mathcal{L}\left(x_{1}, \ldots, x_{n} ; \mu\right)=\sum_{i=1}^{n} f_{i}\left(x_{i}\right)-\mu^{T}\left(\sum_{i} A_{i} x_{i}-b\right)+\frac{\beta}{2}\left\|\sum_{i} A_{i} x_{i}-b\right\|^{2} . \tag{5}
\end{equation*}
$$

The convergence of the direct extension of ADMM to multi-block case had been a long standing open question, until a counter-example was recently given in [5]. More specifically, [5] showed that even for the simplest scenario where the objective function is 0 and the number of blocks is 3 , ADMM can be divergent for a certain choice of $A=\left[A_{1}, A_{2}, A_{3}\right]$ (in fact, there is a positive measure of $A$ such that ADMM can be divergent). There are several proposals to overcome the drawback (see, e.g., $[6,11-|5,17,22|$ ), but they either need to restrict the range of original problems being solved, add additional cost in each step of computation, or limit the stepsize in updating the Lagrangian multipliers. These solutions typically slow down the performance of ADMM for solving most practical problems. One may ask whether a "minimal" modification of cyclic multi-block ADMM (4) can lead to convergence.

One of the simplest modifications of (4) is to add randomness to the update order. Randomness in the update order has been very useful in the analysis of block coordinate gradient descent (BCGD) and stochastic gradient descent (SGD). In particular, the known iteration complexity bounds of randomized BCGD [18] and SAG (Stochastic Average Gradient, a variant of SGD) [20] are much better than the known iteration complexity bounds of their cyclic counterparts BCGD [1] and IAG (Incremental Aggregated Gradient) [2], respectively. ${ }^{1}$ The iteration complexity bounds for randomized algorithms are usually established for independent randomization (sampling with replacement), while in practice, random permutation (sampling without replacement) has been reported to exhibit faster convergence (e.g. [19, 21,23]). However, the theoretical analysis for random permutation seems to be very difficult since the picked blocks/components are not independent across iterations. We have tested both randomly permuted and independently randomized versions of ADMM. Interestingly, independently randomized versions can still be divergent, even for solving linear system of equations, while random permutation can make ADMM converge in all experiments we have conducted.

The main result of this paper is to support the above observation: when the objective function is zero and the constraint is a non-singular square linear system of equations, the expected output of randomly permuted ADMM (RP-ADMM) converges to the unique primal-dual optimal solution. Our contributions are two-fold. First, our result shows that RP-ADMM may serve as a simple solution to resolve the divergence issue of cyclic multi-block ADMM. Since multi-block ADMM is one promising candidate of fast algorithms for large-scale linearly constrained problems, we expect RP-ADMM to be one of the major solvers in big data optimization. Second, our result is one of the first direct analysis of random permutation (sampling without replacement) in optimization algorithms. Our proof framework and techniques will be of independent interest and can be used to analyze random permutation in other optimization algorithms.
We restrict to the simple category of solving linear system of equations, instead of the general convex optimization problems, since the counter-example in [5] belongs to this category and this category seems already difficult to handle for ADMM. The difficulty lies in how to proving the spectral

[^1]radius of the expected update matrix $M$ is less than one. There are two issues: first, there are few mathematical tools to deal with the spectral radius of non-symmetric matrices; second, the entries of $M$ are complicated functions of the entries of $A^{T} A$ (in fact, $n$-th order polynomials). To resolve the first issue, we build a relation between the eigenvalues of $M \in \mathbb{R}^{2 N \times 2 N}$ and the eigenvalues of a symmetric matrix $A Q A^{T} \in \mathbb{R}^{N \times N}$ (see Lemma 1 ), where $Q$ is the expectation of the inverse of a random matrix. To resolve the second issue, we use mathematical induction to implicitly utilize the relation of the entries of $A Q A^{T}$ and $A$. The induction analysis requires several techniques, including a three-level symmetrization technique to construct an induction formula that relates $Q$ to its lower dimensional analogs.

Organization. In Section 2, we present RP-ADMM. Two other versions of randomized ADMM are presented in Section 3 In Section 4 , we present our main results Theorem 1 . Theorem 2 and their proofs. The proofs of the two technical results Lemma 1 and Lemma 2, which are used in the proof of Theorem 2, are provided in the supplement.
Notations. For a matrix $X$, we denote $X(i, j)$ as the $(i, j)$-th entry of $X$, eig $(X)$ as the set of eigenvalues of $X, \rho(X)$ as the spectral radius of $X$ (i.e. the maximum modulus of the eigenvalues of $X$ ), $\|X\|$ as the spectral norm of $X$, and $X^{T}$ as the transpose of $X$. When $X$ is block partitioned, we use $X[i, j]$ to denote the $(i, j)$-th block of $X$. When $X$ is a real symmetric matrix, let $\lambda_{\max }(X)$ and $\lambda_{\min }(X)$ denote the maximum and minimum eigenvalue of $X$ respectively.

## 2 Randomly Permuted ADMM

In this section, we first present RP-ADMM (Randomly Permuted ADMM) for solving the optimization problem (1), then we specialize RP-ADMM for solving a linear system of equations.

Define $\Gamma$ as

$$
\begin{equation*}
\Gamma \triangleq\{\sigma \mid \sigma \text { is a permutation of }\{1, \ldots, n\}\} . \tag{6}
\end{equation*}
$$

At each round, we draw a permutation $\sigma$ of $\{1, \ldots, n\}$ uniformly at random from $\Gamma$, and update the primal variables in the order of the permutation, followed by updating the dual variables in a usual way. Obviously, all primal and dual variables are updated exactly once at each round. See Algorithm 1 for the details of RP-ADMM. Note that with a little abuse of notation, the function $\mathcal{L}\left(x_{\sigma(1)}, x_{\sigma(2)}, \ldots, x_{\sigma(n)} ; \mu\right)$ in this algorithm should be understood as $\mathcal{L}\left(x_{1}, x_{2}, \ldots, x_{n} ; \mu\right)$. For example, when $n=3$ and $\sigma=(231), \mathcal{L}\left(x_{\sigma(1)}, x_{\sigma(2)}, x_{\sigma(3)} ; \mu\right)=\mathcal{L}\left(x_{2}, x_{3}, x_{1} ; \mu\right)$ should be understood as $\mathcal{L}\left(x_{1}, x_{2}, x_{3} ; \mu\right)$.

```
Algorithm \(1 n\)-block Randomly Permuted ADMM (RP-ADMM)
    Initialization: \(x_{i}^{0} \in \mathbb{R}^{d_{i} \times 1}, i=1, \ldots, n ; \mu^{0} \in \mathbb{R}^{N \times 1}\).
    Round \(k(k=0,1,2, \ldots)\) :
    1) Primal update.
        Pick a permutation \(\sigma\) of \(\{1, \ldots, n\}\) uniformly at random.
        For \(i=1, \ldots, n\), compute \(x_{\sigma(i)}^{k+1}\) by
\[
\begin{equation*}
x_{\sigma(i)}^{k+1}=\arg \min _{x_{\sigma(i)} \in \mathcal{X}} \mathcal{X}_{\sigma(i)} \mathcal{L}\left(x_{\sigma(1)}^{k+1}, \ldots, x_{\sigma(i-1)}^{k+1}, x_{\sigma(i)}, x_{\sigma(i+1)}^{k}, \ldots, x_{\sigma(n)}^{k} ; \mu^{k}\right) \tag{7}
\end{equation*}
\]
```

2) Dual update. Update the dual variable by

$$
\begin{equation*}
\mu^{k+1}=\mu^{k}-\beta\left(\sum_{i=1}^{n} A_{i} x_{i}^{k+1}-b\right) . \tag{8}
\end{equation*}
$$

In this paper, we will only consider using Algorithm 1 to solve a square linear system of equations. Consider a special case of (1) where $f_{i}=0, \mathcal{X}_{i}=\mathbb{R}^{d_{i}}, \forall i$ and $N=\sum_{i} d_{i}$ (i.e. the constraint is a square system of equations). Then problem (1) becomes

$$
\begin{array}{rl}
\min _{x \in \mathbb{R}^{N}} & 0  \tag{9}\\
\text { s.t. } & A_{1} x_{1}+\cdots+A_{n} x_{n}=b,
\end{array}
$$

where $A_{i} \in \mathbb{R}^{N \times d_{i}}, x_{i} \in \mathbb{R}^{d_{i} \times 1}, b \in \mathbb{R}^{N \times 1}$. Solving this feasibility problem (with 0 being the objective function) is equivalent to solving a linear system of equations

$$
\begin{equation*}
A x=b \tag{10}
\end{equation*}
$$

where $A=\left[A_{1}, \ldots, A_{n}\right] \in \mathbb{R}^{N \times N}, x=\left[x_{1}^{T}, \ldots, x_{n}^{T}\right]^{T} \in \mathbb{R}^{N \times 1}, b \in \mathbb{R}^{N \times 1}$. Throughout this paper, we assume $A$ is non-singular. Then the unique solution to 10 is $x=A^{-1} b$, and problem (9) has a unique primal-dual optimal solution $(x, \mu)=\left(A^{-1} b, 0\right)$. The augmented Lagrangian function (5) for the optimization problem (9) becomes

$$
\mathcal{L}(x, \mu)=-\mu^{T}(A x-b)+\frac{\beta}{2}\|A x-b\|^{2}
$$

Throughout this paper, we assume $\beta=1$; note that our algorithms and results can be extended to any $\beta>0$ by simply scaling $\mu$.

### 2.1 Example: 3-block ADMM

Before presenting the update equation for general RP-ADMM, we consider a simple case $N=n=$ $3, d_{i}=1, \forall i$ and $\sigma=(123)$, and let $a_{i}=A_{i} \in \mathbb{R}^{3 \times 1}$. The update equations (7) and (8) can be rewritten as

$$
\begin{aligned}
& -a_{1}^{T} \lambda^{k}+a_{1}^{T}\left(a_{1} x_{1}^{k+1}+a_{2} x_{2}^{k}+a_{3} x_{3}^{k}-b\right)=0 \\
& -a_{2}^{T} \lambda^{k}+a_{2}^{T}\left(a_{1} x_{1}^{k+1}+a_{2} x_{2}^{k+1}+a_{3} x_{3}^{k}-b\right)=0 \\
& -a_{3}^{T} \lambda^{k}+a_{3}^{T}\left(a_{1} x_{1}^{k+1}+a_{2} x_{2}^{k+1}+a_{3} x_{3}^{k+1}-b\right)=0 \\
& \left(a_{1} x_{1}^{k+1}+a_{2} x_{2}^{k+1}+a_{3} x_{3}^{k+1}-b\right)+\lambda^{k+1}-\lambda^{k}=0
\end{aligned}
$$

Denote $y^{k}=\left[x_{1}^{k} ; x_{2}^{k} ; x_{3}^{k} ;\left(\lambda^{k}\right)^{T}\right] \in \mathbb{R}^{6 \times 1}$, then the above update equation becomes

$$
\left[\begin{array}{cccc}
a_{1}^{T} a_{1} & 0 & 0 & 0  \tag{11}\\
a_{2}^{T} a_{1} & a_{2}^{T} a_{2} & 0 & 0 \\
a_{3}^{T} a_{1} & a_{3}^{T} a_{2} & a_{3}^{T} a_{3} & 0 \\
a_{1} & a_{2} & a_{3} & I_{3 \times 3}
\end{array}\right] y^{k+1}=\left[\begin{array}{cccc}
0 & -a_{1}^{T} a_{2} & -a_{1}^{T} a_{3} & a_{1}^{T} \\
0 & 0 & -a_{2}^{T} a_{3} & a_{2}^{T} \\
0 & 0 & 0 & a_{3}^{T} \\
0 & 0 & 0 & I_{3 \times 3}
\end{array}\right] y^{k}+\left[\begin{array}{c}
A^{T} b \\
b
\end{array}\right]
$$

Define

$$
L \triangleq\left[\begin{array}{ccc}
a_{1}^{T} a_{1} & 0 & 0  \tag{12}\\
a_{2}^{T} a_{1} & a_{2}^{T} a_{2} & 0 \\
a_{3}^{T} a_{1} & a_{3}^{T} a_{2} & a_{3}^{T} a_{3}
\end{array}\right], \quad R \triangleq\left[\begin{array}{ccc}
0 & -a_{1}^{T} a_{2} & -a_{1}^{T} a_{3} \\
0 & 0 & -a_{2}^{T} a_{3} \\
0 & 0 & 0
\end{array}\right]
$$

The relation between $L$ and $R$ is

$$
L-R=A^{T} A
$$

Define

$$
\bar{L} \triangleq\left[\begin{array}{cc}
L & 0  \tag{13}\\
A & I_{3 \times 3}
\end{array}\right], \quad \bar{R} \triangleq\left[\begin{array}{cc}
R & A^{T} \\
0 & I_{3 \times 3}
\end{array}\right], \quad \bar{b}=\left[\begin{array}{c}
A^{T} b \\
b
\end{array}\right]
$$

then the update equation (11) becomes $\bar{L} y^{k+1}=\bar{R} y^{k}$, i.e.

$$
\begin{equation*}
y^{k+1}=(\bar{L})^{-1} \bar{R} y^{k}+\bar{L}^{-1} \bar{b} \tag{14}
\end{equation*}
$$

As a side remark, reference [5] provides a specific example of $A \in \mathbb{R}^{3 \times 3}$ so that $\rho\left((\bar{L})^{-1} \bar{R}\right)>1$, which implies the divergence of the above iteration if the update order $\sigma=(123)$ is used all the time. This counterexample disproves the convergence of cyclic 3-block ADMM.

### 2.2 General Update Equation of RP-ADMM

In this case, the optimization problem is (9), and the primal update (7) becomes

$$
\begin{equation*}
-A_{\sigma(i)}^{T} \mu^{k}+A_{\sigma(i)}^{T}\left(\sum_{j=1}^{i} A_{\sigma(j)} x_{\sigma(j)}^{k+1}+\sum_{l=i+1}^{n} A_{\sigma(l)} x_{\sigma(l)}^{k}-b\right)=0, i=1, \ldots, n \tag{15}
\end{equation*}
$$

Denote the output of Algorithm 1 after round $(k-1)$ as

$$
y^{k} \triangleq\left[x_{1}^{k} ; \ldots ; x_{n}^{k} ; \mu^{k}\right] \in \mathbb{R}^{2 N \times 1}
$$

Similar to the previous subsection, the update equations of Algorithm 1 for solving (9), i.e. (15) and (8), can be written in the matrix form as (when the permutation is $\sigma$ and $\beta=1$ )

$$
\begin{equation*}
y^{k+1}=\bar{L}_{\sigma}^{-1} \bar{R}_{\sigma} y^{k}+\bar{L}_{\sigma}^{-1} \bar{b}, \tag{16}
\end{equation*}
$$

where $\bar{L}_{\sigma}, \bar{R}_{\sigma}, L_{\sigma}, R_{\sigma}, \bar{b}$ are defined by

$$
\bar{L}_{\sigma} \triangleq\left[\begin{array}{cc}
L_{\sigma} & 0  \tag{17}\\
A & I_{N \times N}
\end{array}\right], \quad \bar{R}_{\sigma} \triangleq\left[\begin{array}{cc}
R_{\sigma} & A^{T} \\
0 & I_{N \times N}
\end{array}\right], \quad \bar{b}=\left[\begin{array}{c}
A^{T} b \\
b
\end{array}\right],
$$

in which $L_{\sigma} \in \mathbb{R}^{N \times N}$ has $n \times n$ blocks and the $(i, j)$-th block is defined as

$$
L_{\sigma}[\sigma(i), \sigma(j)] \triangleq \begin{cases}A_{\sigma(i)}^{T} A_{\sigma(j)} & j \leq i  \tag{18}\\ 0 & j>i\end{cases}
$$

and $R_{\sigma}$ is defined as

$$
R_{\sigma} \triangleq L_{\sigma}-A^{T} A
$$

When $n=3, d_{i}=1, \forall i$ and $\sigma=(123), L_{\sigma}$ defined above is the same as $L$ defined in (12).

## 3 Other Randomized ADMM

In this section, we present two other versions of randomized ADMM which can be divergent according to simulations. The failure of these versions makes us focus on analyzing RP-ADMM in this paper.

In the first algorithm, called primal-dual randomized ADMM (PD-RADMM), the whole dual variable is viewed as the $(n+1)$-th block. In particular, at each iteration, the algorithm draws one index $i$ from $\{1, \ldots, n, n+1\}$, then performs the following update: if $i \leq n$, update the $i$-th block of the primal variable; if $i=n+1$, update the whole dual variable. The details are given in Algorithm 2. We have tested PD-RADMM for the counter-example given in [5], and found that PD-RADMM always diverges (for random initial points).

```
Algorithm 2 Primal-Dual Randomized ADMM (PD-RADMM)
    Iteration \(t(t=0,1,2, \ldots)\) :
        Pick \(i \in\{1, \ldots, n, n+1\}\) uniformly at random;
            If \(1 \leq i \leq n\) :
                        \(x_{i}^{t+\overline{1}}=\arg \min _{x_{i} \in \mathcal{X}_{i}} \mathcal{L}\left(x_{1}^{t}, \ldots, x_{i-1}^{t}, x_{i}, x_{i+1}^{t}, \ldots, x_{n}^{t} ; \mu^{t}\right)\),
                        \(x_{j}^{t+1}=x_{j}^{t}, \forall j \in\{1, \ldots, n\} \backslash\{i\}\),
                        \(\mu^{t+1}=\mu^{t}\).
            Else If \(i=n+1\) :
                        \(\mu^{t+1}=\mu^{t}-\beta\left(\sum_{i=1}^{n} A_{i} x_{i}^{t+1}-b\right)\),
                        \(x_{j}^{t+1}=x_{j}^{t}, \forall j \in\{1, \ldots, n\}\).
            End
```

In the second algorithm, called primal randomized ADMM (P-RADMM), we only perform randomization for the primal variables. In particular, at each round, we first draw $n$ independent random variables $j_{1}, \ldots, j_{n}$ from the uniform distribution of $\{1, \ldots, n\}$ and update $x_{j_{1}}, \ldots, x_{j_{n}}$ sequentially, then update the dual variable in the usual way. The details are given in Algorithm 3. This algorithm looks quite similar to RP-ADMM as they both update $n$ primal blocks at each round; the difference is that RP-ADMM samples without replacement while this algorithm P-RADMM samples with replacement. In other words, RP-ADMM updates each block exactly once at each round, while P-RADMM may update one block more than one times or does not update one block at each round. We have tested P-RADMM in various settings. For the counter-example given in [5], we found that P-RADMM does converge. However, if $n \geq 30$ and $A$ is a Gaussian random matrix (each entry is drawn i.i.d. from $\mathcal{N}(0,1)$ ), then P-RADMM diverges in almost all cases we have tested. This phenomenon is rather strange since for random Gaussian matrices $A$ the cyclic ADMM actually converges (according to simulations). An implication is that randomized versions do not always outperform their deterministic counterparts in terms of convergence.

Since both Algorithm 2 and Algorithm 3 can diverge in certain cases, we will not further study them in this paper. In the rest of the paper, we will focus on RP-ADMM (i.e. Algorithm 1).

```
Algorithm 3 Primal Randomized ADMM (P-RADMM)
    Round \(k(k=0,1,2, \ldots)\) :
    1) Primal update.
        Pick \(l_{1}, \ldots, l_{n}\) independently from the uniform distribution of \(\{1, \ldots, n\}\).
        For \(i=1, \ldots, n\) :
            \(t=k n+i-1\),
            \(x_{l_{i}}^{t+1}=\arg \min _{x_{l_{i}} \in \mathcal{X}_{l_{i}}} \mathcal{L}\left(x_{1}^{t}, \ldots, x_{l_{i}-1}^{t}, x_{l_{i}}, x_{l_{i}+1}^{t}, \ldots, x_{n}^{t} ; \mu^{t}\right)\),
            \(x_{j}^{t+1}=x_{j}^{t}, \forall j \in\{1, \ldots, n\} \backslash\left\{l_{i}\right\}\),
            \(\mu^{t+1}=\mu^{t}\).
```

        End.
    2) Dual update.
        \(\mu^{(k+1) n}=\mu^{k n}-\beta\left(\sum_{i=1}^{n} A_{i} x_{i}^{(k+1) n}-b\right)\).
    
## 4 Main Results

Let $\sigma_{i}$ denote the permutation used in round $i$ of Algorithm 1 , which is a uniform random variable drawn from the set of permutations $\Gamma$. After round $k$, Algorithm 1 generates a random output $y^{k+1}$, which depends on the observed draw of the random variable

$$
\begin{equation*}
\xi_{k}=\left(\sigma_{0}, \sigma_{1}, \ldots, \sigma_{k}\right) \tag{19}
\end{equation*}
$$

We will show that the expected output

$$
\begin{equation*}
\phi^{k}=E_{\xi_{k-1}}\left(y^{k}\right) \tag{20}
\end{equation*}
$$

converges to the primal-dual solution of the problem (9). Note that the expected iterate convergence does not necessarily implies that the iterates converge. However, it strongly indicates that random permutation make a dramatic difference in multi-block ADMM (i.e. ADMM with more than two blocks).
Theorem 1 Assume the coefficient matrix $A=\left[A_{1}, \ldots, A_{n}\right]$ of the constraint in (9] is a nonsingular square matrix. Suppose Algorithm 1 is used to solve problem (9), then the expected output converges to the unique primal-dual optimal solution to (9), i.e.

$$
\left\{\phi^{k}\right\}_{k \rightarrow \infty} \longrightarrow\left[\begin{array}{c}
A^{-1} b  \tag{21}\\
0
\end{array}\right]
$$

Since the update matrix does not depend on previous iterates, we claim (and prove in Section 4.1) that Theorem 1 holds if the expected update matrix has a spectral radius less than 1, i.e. if the following Theorem 2 holds.

Theorem 2 Suppose $A=\left[A_{1}, \ldots, A_{n}\right] \in \mathbb{R}^{N \times N}$ is non-singular, and $\bar{L}_{\sigma}^{-1}, \bar{R}_{\sigma}$ are defined by (17) for any permutation $\sigma$. Define

$$
\begin{equation*}
M \triangleq E_{\sigma}\left(\bar{L}_{\sigma}^{-1} \bar{R}_{\sigma}\right)=\frac{1}{n!} \sum_{\sigma \in \Gamma}\left(\bar{L}_{\sigma}^{-1} \bar{R}_{\sigma}\right) \tag{22}
\end{equation*}
$$

where the expectation is taken over the uniform random distribution over $\Gamma$, the set of permutations of $\{1,2, \ldots, n\}$. Then the spectral radius of $M$ is smaller than 1 , i.e.

$$
\begin{equation*}
\rho(M)<1 \tag{23}
\end{equation*}
$$

Remark 4.1 For the counterexample in [5] where $A=[1,1,1 ; 1,1,2 ; 1,2,2]$, we have $\rho\left(M_{\sigma}\right)>$ 1.02 for any permutation of $(1,2,3)$. Theorem 2 shows that even if each $M_{\sigma}$ is "bad", the average of them is always "good".

Theorem 2 is just a linear algebra result, and can be understood even without knowing the details of the algorithm. However, the proof of Theorem 2 is rather non-trivial and forms the main body of the paper. This proof will be provided in Section 4.2, and the technical results used in this proof will be provided in the supplement.

### 4.1 Proof of Theorem 1

Denote $\sigma_{k}$ as the permutation used in round $k$, and define $\xi_{k}$ as in 19). Rewrite the update equation (16) below (replacing $\sigma$ by $\sigma_{k}$ ):

$$
\begin{equation*}
y^{k+1}=\bar{L}_{\sigma_{k}}^{-1} \bar{R}_{\sigma_{k}} y^{k}+\bar{L}_{\sigma_{k}}^{-1} \bar{b} \tag{24}
\end{equation*}
$$

We first prove (21) for the case $b=0$. By (17) we have $\bar{b}=0$, then (24) is simplified to $y^{k+1}=$ $\bar{L}_{\sigma_{k}}^{-1} \bar{R}_{\sigma_{k}} y^{k}$. Taking the expectation of both sides of this equation in $\xi_{k}$ (see its definition in (19)), and note that $y^{k}$ is independent of $\sigma_{k}$, we get

$$
\phi^{k+1}=E_{\xi_{k}}\left(\bar{L}_{\sigma_{k}}^{-1} \bar{R}_{\sigma_{k}} y^{k}\right)=E_{\sigma_{k}}\left(E_{\xi_{k-1}}\left(\bar{L}_{\sigma_{k}}^{-1} \bar{R}_{\sigma_{k}} y^{k}\right)\right)=E_{\sigma_{k}}\left(\bar{L}_{\sigma_{k}}^{-1} \bar{R}_{\sigma_{k}} \phi^{k}\right)=M \phi^{k}
$$

Since the spectral radius of $M$ is less than 1 by Theorem 2, we have that $\left\{\phi^{k}\right\} \rightarrow 0$, i.e. (21).
We then prove (21) for general $b$. Let $y^{*}=\left[A^{-1} b ; 0\right]$ denote the optimal solution. Then it is easy to verify that

$$
y^{*}=\bar{L}_{\sigma_{k}}^{-1} \bar{R}_{\sigma_{k}} y^{*}+\bar{L}_{\sigma_{k}}^{-1} \bar{b}
$$

for all $\sigma_{k} \in \Gamma$ (i.e. the optimal solution is the fixed point of the update equation for any order). Compute the difference between this equation and (24) and letting $\hat{y}^{k}=y^{k}-y^{*}$, we get $\hat{y}^{k+1}=$ $\bar{L}_{\sigma_{k}}^{-1} \bar{R}_{\sigma_{k}} \hat{y}^{k}$. According to the proof for the case $b=0$, we have $E\left(\hat{y}^{k}\right) \longrightarrow 0$, which implies $E\left(y^{k}\right) \longrightarrow y^{*}$.

### 4.2 Proof of Theorem 2

The difficulty of proving Theorem 2 (bounding the spectral radius of $M$ ) is two-fold. First, $M$ is a non-symmetric matrix, and there are very few tools to bound the spectral radius of a non-symmetric matrix. In fact, spectral radius is neither subadditive nor submultiplicative (see, e.g. [16]). Note that the spectral norm of $M$ can be much larger than 1 (there are examples that $\|M\|>2$ ), thus we cannot bound the spectral radius simply by the spectral norm. Second, although it is possible to explicitly write each entry of $M$ as a function of the entries of $A^{T} A$, these functions are very complicated ( $n$-th order polynomials).
The proof outline of Theorem 2 and the main techniques are described below. In Step 0 , we provide an expression of the expected update matrix $M$. In Step 1, we establish the relationship between the eigenvalues of $M$ and the eigenvalues of a simple symmetric matrix $A Q A^{T}$. As a consequence, the spectral radius of $M$ is smaller than one iff the eigenvalues of $A Q A^{T}$ lie in the region $(0,4 / 3)$. This step partially resolves the first difficulty, i.e. how to deal with the spectral radius of a nonsymmetric matrix. In Step 2, we show that the eigenvalues of $A Q A^{T}$ do lie in $(0,4 / 3)$ using mathematical induction. The induction analysis circumvents the second difficulty, i.e. how to utilize the relation between $M$ and $A$. Note that we will perform induction analysis for $Q A^{T} A$ (with the same eigenvalues as $A Q A^{T}$ ) which is non-symmetric, and we will use several techniques in Step 1 again to transform non-symmetric matrices to symmetric matrices.

Step 0: compute the expression of the expected update matrix $M$. Define

$$
\begin{equation*}
Q \triangleq E_{\sigma}\left(L_{\sigma}^{-1}\right)=\frac{1}{n!} \sum_{\sigma \in \Gamma} L_{\sigma}^{-1} \tag{25}
\end{equation*}
$$

It is easy to prove that $Q$ defined by $[25]$ is symmetric. In fact, note that $L_{\sigma}^{T}=L_{\bar{\sigma}}, \forall \sigma \in \Gamma$, where $\bar{\sigma}$ is a reverse permutation of $\sigma$ satisfying $\bar{\sigma}(i)=\sigma(n+1-i), \forall i$, thus $Q=\frac{1}{n!} \sum_{\sigma} Q_{\sigma}=$ $\left(\frac{1}{n!} \sum_{\sigma} Q_{\bar{\sigma}}\right)^{T}=Q^{T}$, where the last step is because the sum of all $Q_{\bar{\sigma}}$ is the same as the sum of all $Q_{\sigma}$. Denote

$$
M_{\sigma} \triangleq \bar{L}_{\sigma}^{-1} \bar{R}_{\sigma}=\bar{L}_{\sigma}^{-1}\left[\begin{array}{cc}
R_{\sigma} & A^{T}  \tag{26}\\
0 & I
\end{array}\right]
$$

Substituting the expression of $\bar{L}_{\sigma}^{-1}$ into the above relation, and replacing $R_{\sigma}$ by $L_{\sigma}-A^{T} A$, we obtain

$$
M_{\sigma}=\left[\begin{array}{cc}
L_{\sigma}^{-1} & 0  \tag{27}\\
-A L_{\sigma}^{-1} & I
\end{array}\right]\left[\begin{array}{cc}
L_{\sigma}-A^{T} A & A^{T} \\
0 & I
\end{array}\right]=\left[\begin{array}{cc}
I-L_{\sigma}^{-1} A^{T} A & L_{\sigma}^{-1} A^{T} \\
-A+A L_{\sigma}^{-1} A^{T} A & I-A L_{\sigma}^{-1} A^{T}
\end{array}\right] .
$$

Since $M_{\sigma}$ is linear in $L_{\sigma}^{-1}$, we have

$$
\begin{align*}
M=E_{\sigma}\left(M_{\sigma}\right) & =\left[\begin{array}{cc}
I-E_{\sigma}\left(L_{\sigma}^{-1}\right) A^{T} A & E_{\sigma}\left(L_{\sigma}^{-1}\right) A^{T} \\
-A+A E_{\sigma}\left(L_{\sigma}^{-1}\right) A^{T} A & I-A E_{\sigma}\left(L_{\sigma}^{-1}\right) A^{T}
\end{array}\right] \\
& =\left[\begin{array}{cc}
I-Q A^{T} A & Q A^{T} \\
-A+A Q A^{T} A & I-A Q A^{T}
\end{array}\right] . \tag{28}
\end{align*}
$$

Step 1: relate $M$ to a simple symmetric matrix. The main result of Step 1 is given below, and the proof of this result is given in the supplement.

Lemma 1 Suppose $A \in \mathbb{R}^{N \times N}$ is non-singular and $Q \in \mathbb{R}^{N \times N}$ is an arbitrary matrix. Define $M \in \mathbb{R}^{2 N \times 2 N}$ as

$$
M=\left[\begin{array}{cc}
I-Q A^{T} A & Q A^{T}  \tag{29}\\
-A+A Q A^{T} A & I-A Q A^{T}
\end{array}\right] .
$$

Then

$$
\begin{equation*}
\lambda \in \operatorname{eig}(M) \Longleftrightarrow \frac{(1-\lambda)^{2}}{1-2 \lambda} \in \operatorname{eig}\left(Q A^{T} A\right) \tag{30}
\end{equation*}
$$

Furthermore, when $Q$ is symmetric, we have

$$
\begin{equation*}
\rho(M)<1 \Longleftrightarrow \operatorname{eig}\left(Q A^{T} A\right) \subseteq\left(0, \frac{4}{3}\right) \tag{31}
\end{equation*}
$$

Remark: For our problem, the matrix $Q$ as defined by (25) is symmetric (see the argument after equation (25). Lemma 1 implies (31) holds. Note that the first conclusion (30) holds even if $Q$ is non-symmetric.

Step 2: Bound the eigenvalues of $Q A^{T} A$. The main result of Step 2 is summarized in the following Lemma 2. The proof of Lemma 2 is based on an induction formula that relates $Q$ to its lower dimensional analogs, and several techniques used in the proof of Lemma 1, see the supplement for the details of the proof.

Lemma 2 Suppose $A=\left[A_{1}, \ldots, A_{n}\right] \in \mathbb{R}^{N \times N}$ is non-singular. Define $Q$ as

$$
\begin{equation*}
Q \triangleq E_{\sigma}\left(L_{\sigma}^{-1}\right)=\frac{1}{n!} \sum_{\sigma \in \Gamma} L_{\sigma}^{-1} \tag{32}
\end{equation*}
$$

in which $L_{\sigma}$ is defined by (18) and $\Gamma$ is defined by (6). Then all eigenvalues of $Q A^{T} A$ lie in $(0,4 / 3)$, i.e.

$$
\begin{equation*}
\operatorname{eig}\left(Q A^{T} A\right) \subseteq\left(0, \frac{4}{3}\right) \tag{33}
\end{equation*}
$$

Theorem 2 follows immediately from Lemma 1 and Lemma 2

## 5 Conclusion

In this paper, we propose randomly permuted ADMM (RP-ADMM) and prove the expected convergence of RP-ADMM for solving a non-singular square system of equations. Multi-block ADMM is one promising candidate for solving large-scale linearly constrained problems in big data applications, but its cyclic version is known to be possibly divergent. Our result shows that RP-ADMM may serve as a simple solution to resolve the divergence issue of cyclic multi-block ADMM. One interesting aspect is that while it is possible that every single permutation leads to a "bad" update matrix, averaging these permutations always leads to a "good" update matrix. Our result is also one of the first direct analysis of random permutation (sampling without replacement) in optimization algorithms, though independent randomization (sampling with replacement) has been extensively studied for BCD and SGD. It is not hard to extend our result to non-square (including tall and wide) system of equations. Future directions include extending our result to general convex problems and proving the convergence of RP-ADMM with high probability.

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## Supplemental Materials

## 6 Proof of Lemma 1

The proof of Lemma11relies on two simple techniques. The first technique, as elaborated in the Step 1 below, is to factorize $M$ and rearrange the factors. The second technique, as elaborated in the Step 2 below, is to reduce the dimension by eliminating a variable from the eigenvalue equation.

Step 1: Factorizing $M$ and rearranging the order of multiplication. The following observation is crucial: the matrix $M$ defined by (29) can be factorized as

$$
M=\left[\begin{array}{cc}
I & 0 \\
-A & I
\end{array}\right]\left[\begin{array}{cc}
Q A^{T} & I \\
I & A
\end{array}\right]\left[\begin{array}{cc}
-A & I \\
I & 0
\end{array}\right] .
$$

Switching the order of the products by moving the first component to the last, we get a new matrix

$$
M^{\prime} \triangleq\left[\begin{array}{cc}
Q A^{T} & I  \tag{34}\\
I & A
\end{array}\right]\left[\begin{array}{cc}
-A & I \\
I & 0
\end{array}\right]\left[\begin{array}{cc}
I & 0 \\
-A & I
\end{array}\right]=\left[\begin{array}{cc}
Q A^{T} & I \\
I & A
\end{array}\right]\left[\begin{array}{cc}
-2 A & I \\
I & 0
\end{array}\right]=\left[\begin{array}{cc}
I-2 Q A^{T} A & Q A^{T} \\
-A & I
\end{array}\right] .
$$

Note that eig $(X Y)=\operatorname{eig}(Y X)$ for any two square matrices, thus

$$
\operatorname{eig}(M)=\operatorname{eig}\left(M^{\prime}\right)
$$

To prove (30), we only need to prove

$$
\begin{equation*}
\lambda \in \operatorname{eig}\left(M^{\prime}\right) \Longleftrightarrow \frac{(1-\lambda)^{2}}{1-2 \lambda} \in \operatorname{eig}\left(Q A^{T} A\right) . \tag{35}
\end{equation*}
$$

Step 2: Relate the eigenvalues of $M^{\prime}$ to the eigenvalues of $Q A^{T} A$, i.e. prove 35). This step is simple as we only use the definition of eigenvalues. However, note that, without Step 1, just applying the definition of eigenvalues of the original matrix $M$ may not lead to a simple relationship as (35).
We first prove one direction of (30):

$$
\begin{equation*}
\lambda \in \operatorname{eig}\left(M^{\prime}\right) \Longrightarrow \frac{(1-\lambda)^{2}}{1-2 \lambda} \in \operatorname{eig}\left(Q A^{T} A\right) \tag{36}
\end{equation*}
$$

Suppose $v \in \mathbb{C}^{2 N \times 1} \backslash\{0\}$ is an eigenvector of $M^{\prime}$ corresponding to the eigenvalue $\lambda$, i.e.

$$
M^{\prime} v=\lambda v .
$$

Partition $v$ as $v=\left[\begin{array}{l}v_{1} \\ v_{0}\end{array}\right]$, where $v_{1}, v_{0} \in \mathbb{C}^{N \times 1}$. Using the expression of $M^{\prime}$ in (34), we can write the above equation as

$$
\left[\begin{array}{cc}
I-2 Q A^{T} A & Q A^{T} \\
-A & I
\end{array}\right]\left[\begin{array}{l}
v_{1} \\
v_{0}
\end{array}\right]=\lambda\left[\begin{array}{l}
v_{1} \\
v_{0}
\end{array}\right],
$$

which implies

$$
\begin{align*}
\left(I-2 Q A^{T} A\right) v_{1}+Q A^{T} v_{0} & =\lambda v_{1}  \tag{37a}\\
-A v_{1}+v_{0} & =\lambda v_{0} \tag{37b}
\end{align*}
$$

We claim that (36) holds when $v_{1}=0$. In fact, in this case we must have $v_{0} \neq 0$ (otherwise $v=0$ cannot be an eigenvector). By (37b) we have $\lambda v_{0}=v_{0}$, thus $\lambda=1$. By 37a we have $0=Q A^{T} v_{0}=Q A^{T} A\left(A^{-1} v_{0}\right)$, which implies $\frac{(1-\lambda)^{2}}{1-2 \lambda}=0 \in \operatorname{eig}\left(Q A^{T} A\right)$, therefore (36) holds in this case.

We then prove (36) for the case

$$
\begin{equation*}
v_{1} \neq 0 \tag{38}
\end{equation*}
$$

The equation (37b) implies $(1-\lambda) v_{0}=A v_{1}$. Multiplying both sides of 37a by $(1-\lambda)$ and invoking this equation, we get

$$
(1-\lambda)\left(I-2 Q A^{T} A\right) v_{1}+Q A^{T} A v_{1}=(1-\lambda) \lambda v_{1}
$$

This relation can be simplified to

$$
\begin{equation*}
(1-2 \lambda) Q A^{T} A v_{1}=(1-\lambda)^{2} v_{1} . \tag{39}
\end{equation*}
$$

We must have $\lambda \neq \frac{1}{2}$; otherwise, the above relation implies $v_{1}=0$, which contradicts 38). Then 39) becomes

$$
\begin{equation*}
Q A^{T} A v_{1}=\frac{(1-\lambda)^{2}}{1-2 \lambda} v_{1} \tag{40}
\end{equation*}
$$

Therefore, $\frac{(1-\lambda)^{2}}{1-2 \lambda}$ is an eigenvalue of $Q A^{T} A$, with the corresponding eigenvector $v_{1} \neq 0$, which finishes the proof of 36.
The other direction $\square^{2}$

$$
\begin{equation*}
\lambda \in \operatorname{eig}(M) \Longleftarrow \frac{(1-\lambda)^{2}}{1-2 \lambda} \in \operatorname{eig}\left(Q A^{T} A\right) \tag{41}
\end{equation*}
$$

is easy to prove. Suppose $\frac{(1-\lambda)^{2}}{1-2 \lambda} \in \operatorname{eig}\left(Q A^{T} A\right)$. We consider two cases.
Case 1: $\frac{(1-\lambda)^{2}}{1-2 \lambda}=0$. In this case $\lambda=1$. Since $0=\frac{(1-\lambda)^{2}}{1-2 \lambda} \in \operatorname{eig}\left(Q A^{T} A\right)$, there exists $v_{0} \in \mathbb{C}^{N} \backslash\{0\}$ such that $Q A^{T} A v_{0}=0$ and Let $v_{1}=(0, \ldots, 0)^{T} \in \mathbb{C}^{N \times 1}$, then $v_{0}, v_{1}$ and $\lambda=1$ satisfy 37). Thus $v=\left[\begin{array}{l}v_{1} \\ v_{0}\end{array}\right] \in \mathbb{C}^{2 N} \backslash\{0\}$ satisfies $M v=\lambda v$, which implies $\lambda=1 \in \operatorname{eig}(M)$.

Case 2: $\frac{(1-\lambda)^{2}}{1-2 \lambda} \neq 0$, then $\lambda \neq 1$. Let $v_{1}$ be the eigenvector corresponding to $\frac{(1-\lambda)^{2}}{1-2 \lambda}$ (i.e. pick $v_{1}$ that satisfies (40), and define $v_{0}=v_{1} /(1-\lambda)$. It is easy to verify that $v=\left[\begin{array}{l}v_{1} \\ v_{0}\end{array}\right]$ satisfies $M v=\lambda v$, which implies $\lambda \in \operatorname{eig}(M)$.
Step 3: When $Q$ is symmetric, prove 31 by simple algebraic computation.
Since $Q$ is symmetric, we know that $\operatorname{eig}\left(Q A^{T} A\right)=\operatorname{eig}\left(A Q A^{T}\right) \subseteq \mathbb{R}$. Suppose $\tau \in \mathbb{R}$ is an eigenvalue of $Q A^{T} A$, then any $\lambda$ satisfying $\frac{(1-\lambda)^{2}}{1-2 \lambda}=\tau$ is an eigenvalue of $M$. This relation can be rewritten as $\lambda^{2}+2(\tau-$ 1) $\lambda+(1-\tau)=0$, which, as a real-coefficient quadratic equation in $\lambda$, has two roots

$$
\begin{equation*}
\lambda_{1}=1-\tau+\sqrt{\tau(\tau-1)}, \quad \lambda_{2}=1-\tau-\sqrt{\tau(\tau-1)} . \tag{42}
\end{equation*}
$$

Note that when $\tau(\tau-1)<0$, the expression $\sqrt{\tau(\tau-1)}$ denotes a complex number $i \sqrt{\tau(1-\tau)}$, where $i$ is the imaginary unit. To prove (31), we only need to prove

$$
\begin{equation*}
\max \left\{\left|\lambda_{1}\right|,\left|\lambda_{2}\right|\right\}<1 \Longleftrightarrow 0<\tau<\frac{4}{3} \tag{43}
\end{equation*}
$$

Consider three cases.
Case 1: $\tau<0$. Then $\tau(\tau-1)=|\tau|(|\tau|+1)>0$. In this case, $\lambda_{1}=1+|\tau|+\sqrt{|\tau|(|\tau|+1)}>1$.
Case 2: $0<\tau<1$. Then $\tau(\tau-1)<0$, and (42) can be rewritten as

$$
\lambda_{1,2}=1-\tau \pm i \sqrt{\tau(1-\tau)}
$$

which implies $\left|\lambda_{1}\right|=\left|\lambda_{2}\right|=\sqrt{(1-\tau)^{2}+\tau(1-\tau)}=\sqrt{1-\tau}<1$.
Case 3: $\tau>1$. Then $\tau(\tau-1)>0$. According to (42), it is easy to verify $\lambda_{1}>0>\lambda_{2}$ and

$$
\left|\lambda_{2}\right|=\tau-1+\sqrt{\tau(\tau-1)}>1-\tau+\sqrt{\tau(\tau-1)}=\left|\lambda_{1}\right| .
$$

Then we have

$$
\max \left\{\left|\lambda_{1}\right|,\left|\lambda_{2}\right|\right\}<1 \Longleftrightarrow\left|\lambda_{2}\right|=\tau-1+\sqrt{\tau(\tau-1)}<1 \Longleftrightarrow 1<\tau<\frac{4}{3} .
$$

Combining the conclusions of the three cases immediately leads to (43).

## 7 Proof of Lemma 2 for the case $d_{i}=1, \forall i$

In this section, we prove Lemma 2 for the case $d_{i}=1, \forall i$. The proof for general $d_{i}$ 's is quite similar (but not exactly the same), and will be given in Section 8

### 7.1 Proof Overview

We will use mathematical induction to prove Lemma 2 and the reason of doing so is the following. A major difficulty of proving Lemma 2 is that each entry of $Q$ is a complicated function (in fact, $n$-th order polynomial) of the entries of $A^{T} A$. To circumvent this difficulty, we will implicitly exploit the property of $Q$ by an induction analysis on $n$, the number of blocks.

[^2]The difficulty of using induction to prove Lemma 2 is two-fold. First, it is not obvious how $Q$ is related to an analogous matrix in a lower dimension. Second, the simulations show that $\left\|Q A^{T} A\right\|<\frac{4}{3} \ll\|Q\|\left\|A^{T} A\right\|$, thus we have to bound the eigenvalues of the product $Q A^{T} A$, instead of the eigenvalues of $Q$. Even if we know the relationship between $Q$ and a lower-dimensional matrix $\hat{Q}$, it is not obvious how eig $\left(Q A^{T} A\right.$ ) and $\operatorname{eig}\left(\hat{Q} \hat{A}^{T} \hat{A}\right)$ are related, where $\hat{A}$ is a lower-dimensional analog of $A$.
The proof outline of Lemma 2 and the main techniques are described below. In Step 1, we prove an induction formula in Proposition 1 , which states that $Q$ can be decomposed as the sum of $n$ symmetric matrices, where each symmetric matrix contains an $(n-1) \times(n-1)$ sub-matrix $\hat{Q}_{k}$ that is analogous to $Q$. In other words, we relate $Q$ to $n$ analogous matrices $\hat{Q}_{k}, k=1, \ldots, n$ in a lower dimension. To prove the induction formula, we use a three-level symmetrization technique. This induction formula resolves the first difficulty. In Step 2, we prove the induction step, i.e. under the induction hypothesis that $\operatorname{eig}\left(\hat{Q}_{k} \hat{A}_{k}^{T} \hat{A}_{k}\right) \subseteq\left(0, \frac{4}{3}\right), k=1, \ldots, n$, where $\hat{A}_{k}$ is a certain sub-matrix of $A$, the desired result $\operatorname{eig}\left(Q A^{T} A\right) \subseteq\left(0, \frac{4}{3}\right)$ holds. To build the relation between $\operatorname{eig}\left(Q A^{T} A\right)$ and $\operatorname{eig}\left(\hat{Q}_{k} \hat{A}_{k}^{T} \hat{A}_{k}\right)$, we will apply the two simple techniques used in the proof of Lemma 1 factorize and rearrange, and reduce the dimension by eliminating a variable from the eigenvalue equation. Nevertheless, the subsequent analysis is more complicated than the proof of Lemma 1

### 7.2 Proof of Lemma 2 for $d_{i}=1, \forall i$ and Two Propositions

Without loss of generality, we can assume $\left\|a_{i}\right\|^{2}=1, \forall i$ (see the first paragraph of Section 8 for an explanation).

We use mathematical induction to prove Lemma 2 for the $n$-coordinate case. For the basis of the induction ( $n=1$ ), Lemma 2 holds since $Q A^{T} A=1$. Assume Lemma 2 holds for $n-1$, we will prove Lemma 2 for $n$.

### 7.2.1 Step 1: Induction formula for general $n$

Since $d_{i}=1, \forall i$, we have $N=\sum_{i} d_{i}=n$. Denote $a_{i} \triangleq A_{i} \in \mathbb{R}^{n \times 1}(i=1, \ldots, n)$. Denote $[n] \triangleq$ $\{1, \ldots, n\}$. For any $k \in[n]$, define

$$
\begin{equation*}
\Gamma_{k} \triangleq\left\{\sigma^{\prime} \mid \sigma^{\prime} \text { is a permutation of }[n] \backslash\{k\}\right\} . \tag{44}
\end{equation*}
$$

For any $\sigma^{\prime} \in \Gamma_{k}$, similar to (??), we can define $L_{\sigma^{\prime}}, Q_{k} \in \mathbb{R}^{(n-1) \times(n-1)}$ as

$$
\begin{gather*}
L_{\sigma^{\prime}}\left(\sigma^{\prime}(i), \sigma^{\prime}(j)\right) \triangleq \begin{cases}a_{\sigma^{\prime}(i)}^{T} a_{\sigma^{\prime}(j)} & i \geq j, \\
0 & i<j,\end{cases}  \tag{45}\\
\hat{Q}_{k} \triangleq \frac{1}{\left|\Gamma_{k}\right|} \sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\sigma^{\prime}}^{-1}, k=1, \ldots, n . \tag{46}
\end{gather*}
$$

Note that $L_{\sigma^{\prime}}$ and $\hat{Q}_{k}$ are lower-dimensional analogs of $L_{\sigma}$ and $Q$ respectively.
Define $w_{k}$ as the $k$-th column of $A^{T} A$ excluding the entry $a_{k}^{T} a_{k}$, i.e.

$$
\begin{equation*}
w_{k} \triangleq\left[a_{1}, \ldots, a_{k-1}, a_{k+1}, \ldots, a_{n}\right]^{T} a_{k} \in \mathbb{R}^{(n-1) \times 1} \tag{47}
\end{equation*}
$$

Define permutation matrices $S_{1}, \ldots, S_{n} \in \mathbb{R}^{n \times n}$ as follows:

$$
\begin{equation*}
S_{k}(i, i)=1, i=1, \ldots, k-1 ; \quad S_{k}(k+1, k)=\cdots=S_{k}(n, n-1)=1 ; \quad S_{k}(k, n)=1, \tag{48}
\end{equation*}
$$

and all other entries of $S_{k}$ are zero. $S_{k}$ is called a permutation matrix since it corresponds to a permutation $(1, \ldots, k-1, k+1, \ldots, n, k)$; in fact, $(1,2, \ldots, n) S_{k}=(1, \ldots, k-1, k+1, \ldots, n, k)$. Replacing $1,2, \ldots, n$ by column vectors $b_{1}, \ldots, b_{n}$, we get $\left(b_{1}, b_{2}, \ldots, b_{n}\right) S_{k}=\left(b_{1}, \ldots, b_{k-1}, b_{k+1}, \ldots, b_{n}, b_{k}\right)$. This relation can be interpreted as the following column-moving property of $S_{k}$ : right-multiply a matrix by $S_{k}$ will move the $k$-th column to the end (i.e. in the new matrix it becomes the last column). Similarly, $S_{k}^{T}$ has the following row-moving property: left-multiply a matrix by $S_{k}^{T}$ will move the $k$-th row to the end. Note that $S_{n}$ is the identity matrix. Another property is

$$
\begin{equation*}
S_{k}^{T}=S_{k}^{-1} \tag{49}
\end{equation*}
$$

We give an example to illustrate the expressions of $S_{k}$. When $n=3, S_{1}, S_{2}, S_{3} \in \mathbb{R}^{3 \times 3}$ defined in 48) can be explicitly written as

$$
S_{1} \triangleq\left[\begin{array}{lll}
0 & 0 & 1 \\
1 & 0 & 0 \\
0 & 1 & 0
\end{array}\right], S_{2} \triangleq\left[\begin{array}{lll}
1 & 0 & 0 \\
0 & 0 & 1 \\
0 & 1 & 0
\end{array}\right], S_{3}=I_{3 \times 3}
$$

The column-moving property means $\left[b_{1}, b_{2}, b_{3}\right] S_{1}=\left[b_{2}, b_{3}, b_{1}\right]$, and $\left[b_{1}, b_{2}, b_{3}\right] S_{2}=\left[b_{1}, b_{3}, b_{2}\right]$. Similarly, the row-moving property means $S_{1}^{T}\left[\begin{array}{l}b_{1}^{T} \\ b_{2}^{T} \\ b_{3}^{T}\end{array}\right]=\left[\begin{array}{l}b_{2}^{T} \\ b_{3}^{T} \\ b_{1}^{T}\end{array}\right]$ and $S_{2}^{T}\left[\begin{array}{l}b_{1}^{T} \\ b_{2}^{T} \\ b_{3}^{T}\end{array}\right]=\left[\begin{array}{l}b_{1}^{T} \\ b_{3}^{T} \\ b_{2}^{T}\end{array}\right]$.
With these definitions, we are ready to present the induction formula, which builds a relation between $Q$ and its lower-dimensional analogs $\hat{Q}_{k}, k=1, \ldots, n$.

Proposition 1 The matrix $Q=\frac{1}{|\Gamma|} \sum_{\sigma \in \Gamma} L_{\sigma}^{-1}$, where $L_{\sigma}, \Gamma$ are defined by (??) and (6) respectively, can be decomposed as follows:

$$
\begin{equation*}
Q=\frac{1}{n} \sum_{k=1}^{n} S_{k} Q_{k} S_{k}^{T} \tag{50}
\end{equation*}
$$

where

$$
Q_{k} \triangleq\left[\begin{array}{cc}
\hat{Q}_{k} & -\frac{1}{2} \hat{Q}_{k} w_{k}  \tag{51}\\
-\frac{1}{2} w_{k}^{T} \hat{Q}_{k} & 1
\end{array}\right],
$$

in which $\hat{Q}_{k}$ is defined by 46.

The proof of Proposition 1 for the case $n=3$ will be given in Section 7.3 We relegate the proof of Proposition 1 for general $n$ to Appendix 9

### 7.2.2 Step 2: bounding eigenvalues of each $Q_{k}$

According to 50, we have

$$
A Q A^{T}=\frac{1}{n} \sum_{k=1}^{n} A S_{k} Q_{k} S_{k}^{T} A^{T}
$$

Note that $\hat{Q}_{k}$ defined by (46) is symmetric, thus $Q_{k}$ defined by (51) is symmetric, which implies that each $A S_{k} Q_{k} S_{k}^{T} A^{T}$ is a symmetric matrix. From the above relation, we have

$$
\begin{equation*}
\frac{1}{n} \sum_{k=1}^{n} \lambda_{\min }\left(A S_{k} Q_{k} S_{k}^{T} A^{T}\right) \leq \lambda_{\min }\left(A Q A^{T}\right) \leq \lambda_{\max }\left(A Q A^{T}\right) \leq \frac{1}{n} \sum_{k=1}^{n} \lambda_{\max }\left(A S_{k} Q_{k} S_{k}^{T} A^{T}\right) \tag{52}
\end{equation*}
$$

Therefore, to prove $\operatorname{eig}\left(A Q A^{T}\right) \subseteq(0,4 / 3)$, we only need to prove for any $k=1, \ldots, n$,

$$
\begin{equation*}
\operatorname{eig}\left(A S_{k} Q_{k} S_{k}^{T} A^{T}\right)=\operatorname{eig}\left(Q_{k} S_{k}^{T} A^{T} A S_{k}\right) \subseteq(0,4 / 3) \tag{53}
\end{equation*}
$$

By the column moving property of $S_{k}$, we have

$$
\begin{equation*}
\bar{A}_{k} \triangleq A S_{k}=\left[\hat{A}_{k}, a_{k}\right], \tag{54}
\end{equation*}
$$

where $\hat{A}_{k} \triangleq\left[a_{1}, \ldots, a_{k-1}, a_{k+1}, \ldots, a_{n}\right]$. Note that $\hat{Q}_{k}$ only depends on the entries of $\hat{A}_{k}^{T} \hat{A}_{k} \in$ $\mathbb{R}^{(n-1) \times(n-1)}$, thus by the induction hypothesis, we have

$$
\begin{equation*}
\operatorname{eig}\left(\hat{Q}_{k} \hat{A}_{k}^{T} \hat{A}_{k}\right) \subseteq(0,4 / 3) \tag{55}
\end{equation*}
$$

We claim that (53) follows from the induction hypothesis (55) and the expressions (54) and (51). In fact, the following Proposition 2 directly proves (53) for $k=n$. If we replace $A, \hat{A}_{n}, a_{n}, \hat{Q}_{n}, Q_{n}$ by $\bar{A}_{k}, \hat{A_{k}}, a_{k}, \hat{Q}_{k}, Q_{k}$ respectively in Proposition 2 we will obtain 53) for any $k$. As mentioned earlier, the desired result $\operatorname{eig}\left(A Q A^{T}\right) \subseteq(0,4 / 3)$ follows immediately from (53) and (52).

Proposition 2 Suppose $A=\left[\hat{A}_{n}, a_{n}\right] \in \mathbb{R}^{n \times n}$ is non-singular, where $\hat{A}_{n} \in \mathbb{R}^{n \times(n-1)}$ and $a_{n} \in \mathbb{R}^{(n-1) \times 1}$ satisfies $\left\|a_{n}\right\|=1$. Suppose $\hat{Q}_{n} \in \mathbb{R}^{(n-1) \times(n-1)}$ is a symmetric matrix which satisfies eig $\left(\hat{Q}_{n} \hat{A}_{n}^{T} \hat{A}_{n}\right) \subseteq$ $\left(0, \frac{4}{3}\right)$. Define

$$
w_{n} \triangleq \hat{A}_{n}^{T} a_{n}, \quad Q_{n} \triangleq\left[\begin{array}{cc}
\hat{Q}_{n} & -\frac{1}{2} \hat{Q}_{n} w_{n}  \tag{56}\\
-\frac{1}{2} w_{n}^{T} \hat{Q}_{n} & 1
\end{array}\right] .
$$

Then eig $\left(Q_{n} A^{T} A\right) \subseteq\left(0, \frac{4}{3}\right)$.
The proof of Proposition 2 is given in Section 7.4

### 7.3 Proof of Proposition 1 for $n=3$

In this subsection, we prove the induction formula for $n=3$. Before the formal proof, we briefly describe the ideas of contructing this induction formula. The key idea is symmetrization: we start from an obvious relation between $L_{\sigma}^{-1}$ and its lower-dimensional analog, and by three levels of symmetrization we can obtain a relation between $Q$ and its lower-dimensional analogs.

Define

$$
\begin{align*}
w_{i j} & =a_{i}^{T} a_{j}, \forall i, j, \\
w_{1}=\left[w_{12}, w_{13}\right]^{T}, w_{2} & =\left[w_{21}, w_{23}\right]^{T}, w_{3}=\left[w_{31}, w_{32}\right]^{T} . \tag{57}
\end{align*}
$$

For $\sigma=(123)$, the expressions of $L_{\sigma}$ and $L_{\sigma}^{-1}$ are

$$
L_{(123)}=\left[\begin{array}{ccc}
1 & 0 & 0  \tag{58}\\
w_{21} & 1 & 0 \\
w_{31} & w_{32} & 1
\end{array}\right], \quad L_{(123)}^{-1}=\left[\begin{array}{ccc}
1 & 0 & 0 \\
-w_{21} & 1 & 0 \\
-w_{31}+w_{21} w_{32} & -w_{32} & 1
\end{array}\right] .
$$

Note, however, that the following expressions of $L_{\sigma}$ and $L_{\sigma}^{-1}$ are more useful:

$$
L_{(123)}=\left[\begin{array}{cc}
L_{(12)} & 0 \\
w_{3}^{T} & 1
\end{array}\right]
$$

and

$$
L_{(123)}^{-1}=\left[\begin{array}{cc}
L_{(12)}^{-1} & 0  \tag{59}\\
-w_{3}^{T} L_{(12)}^{-1} & 1
\end{array}\right] .
$$

The above equation provides a relation between $L_{(123)}^{-1}$ and an analogous matrix $L_{12}^{-1}$ in a lower dimension. Such a kind of relation also exists between any $L_{\sigma}^{-1}$ and $L_{\sigma^{\prime}}^{-1}$ where $\sigma^{\prime}$ is a sub-permutation of $\sigma$. Here, we say $\sigma^{\prime} \in \Gamma_{k}$ is a sub-permutation of $\sigma \in \Gamma$ if $\sigma^{\prime}(j)=\sigma(j), \forall j \in \Gamma_{k}$. For example, (134) and (123) are both sub-permutations of (1234).
A natural question is: given the relation between $L_{\sigma}$ and its lower dimensional analogs, how to build a relation between $Q$ and its lower dimensional counterparts? To answer this question, the following intuition is crucial: since $Q=E_{\sigma}\left(L_{\sigma}^{-1}\right)$ is a symmetrization of $L_{\sigma}^{-1}$, we should symmetrize RHS (Right-Hand-Side) of (59). There are three levels of "asymmetry" in the RHS of (59): i) $L_{(12)}^{-1}$ is non-symmetric; ii) the off-diagonal blocks are not transpose to each other; iii) in this block partition of $L_{\sigma}$ the 1st and 2nd row/columns are grouped together, so this expression is not symmetric with respect to the permutation of $\{1,2,3\}$. Let us briefly explain below how to build the three levels of symmetry.

The first level of symmetry is built by the matrix $\hat{Q}_{k}$. For example,

$$
\hat{Q}_{3}=\frac{1}{2}\left(L_{(12)}^{-1}+L_{(21)}^{-1}\right)=\left[\begin{array}{cc}
1 & -\frac{1}{2} w_{12}  \tag{60}\\
-\frac{1}{2} w_{12} & 1
\end{array}\right]
$$

is a symmetrization of $L_{(12)}^{-1}$ and forms the first-level symmetrization of the RHS of 59) (more details are given later). The second level of symmetry is built by the matrix $Q_{k}$. For example, $Q_{3}=\left[\begin{array}{cc}\hat{Q}_{3} & -\frac{1}{2} \hat{Q}_{3} w_{3} \\ -\frac{1}{2} w_{3}^{T} \hat{Q}_{3} & 1\end{array}\right]$ is the symmetrization of $\left[\begin{array}{cc}\hat{Q}_{3} & 0 \\ -w_{3}^{T} \hat{Q}_{3} & 1\end{array}\right]$, thus forming the second level of symmetrization for the RHS of (59). The third level of symmetry is built by averaging the three matrices $Q_{1}, Q_{2}, Q_{3}$ (up to permutation of rows/columns), as shown by the induction formula (50)

$$
\begin{equation*}
Q=\frac{1}{3}\left(S_{1} Q_{1} S_{1}^{T}+S_{2} Q_{2} S_{2}^{T}+Q_{3}\right) \tag{61}
\end{equation*}
$$

Below, we prove the induction formula $\sqrt{61]}$ in a rigous way.
Proof of (61): As the first level symmetrization, we prove

$$
\frac{1}{2}\left(L_{(123)}^{-1}+L_{(213)}^{-1}\right)=\left[\begin{array}{cc}
\hat{Q}_{3} & 0  \tag{62}\\
-w_{3}^{T} \hat{Q}_{3} & 1
\end{array}\right]
$$

Recall that $L_{(123)}=\left[\begin{array}{cc}L_{(12)} & 0 \\ w_{3}^{T} & 1\end{array}\right]$ implies $L_{(123)}^{-1}=\left[\begin{array}{cc}L_{(12)}^{-1} & 0 \\ -w_{3}^{T} L_{(12)}^{-1} & 1\end{array}\right]$. Similarly, $L_{(213)}=\left[\begin{array}{cc}L_{(21)} & 0 \\ w_{3}^{T} & 1\end{array}\right]$ implies $L_{(213)}^{-1}=\left[\begin{array}{cc}L_{(21)}^{-1} & 0 \\ -w_{3}^{T} L_{(21)}^{-1} & 1\end{array}\right]$. Summing up these two relations and invoking the definition of $\hat{Q}_{3}$ in (60), we obtain 62).

As the second level symmetrization, we prove

$$
\begin{equation*}
Q_{3}=\frac{1}{4}\left(L_{(123)}^{-1}+L_{(213)}^{-1}+L_{(321)}^{-1}+L_{(312)}^{-1}\right) . \tag{63}
\end{equation*}
$$

Note that the common feature of the four permutations (123), (213), (321), (312) is: 1 and 2 are adjacent in these permutations. By the definition of $L_{\sigma}$ in (??), we have $L_{(321)}=L_{(123)}^{T}, L_{(312)}=L_{(213)}^{T}$, thus $L_{(321)}^{-1}=L_{(123)}^{-T}, L_{(312)}^{-1}=L_{(213)}^{-T}$. Taking the transpose over both sides of 62), we obtain

$$
\frac{1}{2}\left(L_{(321)}^{-1}+L_{(312)}^{-1}\right)=\left[\begin{array}{cc}
\hat{Q}_{3} & -\hat{Q}_{3} w_{3}  \tag{64}\\
0 & 1
\end{array}\right] .
$$

Combining (62) and (64), and using the definition of $Q_{3}$ in (51), we obtain (63).
Using a similar argument, we can prove

$$
\begin{equation*}
S_{1} Q_{1} S_{1}^{T}=\frac{1}{4}\left(L_{(123)}^{-1}+L_{(132)}^{-1}+L_{(231)}^{-1}+L_{(321)}^{-1}\right) \tag{65}
\end{equation*}
$$

Again, the common feature of the four permutations (123), (132), (231), (321) is: 2 and 3 are adjacent in these permutations. The proof of (65) is almost the same as the proof of 63), except the extra step to move rows and columns. Similarly, we can prove

$$
\begin{equation*}
S_{2} Q_{2} S_{2}^{T}=\frac{1}{4}\left(L_{(132)}^{-1}+L_{(312)}^{-1}+L_{(231)}^{-1}+L_{(213)}^{-1}\right) . \tag{66}
\end{equation*}
$$

As the third level symmetrization, combining (63), (65) and (66), and invoking the definition of $Q$ in (25), we obtain (61).

### 7.4 Proof of Proposition 2

For simplicity, throughout this proof, we denote

$$
w \triangleq w_{n}, \hat{Q} \triangleq \hat{Q}_{n}, \hat{A} \triangleq \hat{A}_{n}
$$

We claim that

$$
\begin{equation*}
0 \leq \theta \triangleq w^{T} \hat{Q} w<\frac{4}{3} \tag{67}
\end{equation*}
$$

In fact, by the definition $w=\hat{A}^{T} a_{n}$ we have $\theta=a_{n}^{T} \hat{A} \hat{Q} \hat{A}^{T} a_{n} \leq \rho\left(\hat{A} \hat{Q} \hat{A}^{T}\right)\left\|a_{n}\right\|^{2}=\rho\left(\hat{A} \hat{Q} \hat{A}^{T}\right)<\frac{4}{3}$, which proves the last inequality of (67). According to the assumption, $\operatorname{eig}\left(\hat{Q} \hat{A}^{T} \hat{A}\right) \subseteq(0,4 / 3) \subseteq(0, \infty)$ and $\hat{A}$ is non-singular, thus $\hat{Q} \succ 0$. Then we have $\theta=w^{T} \hat{Q} w \geq 0$, which proves the first inequality of 67).
We apply a trick that we have previously used: factorize $Q_{n}$ and change the order of multiplication. To be specific, $Q_{n}$ defined in (56) can be factorized as

$$
Q_{n}=\left[\begin{array}{cc}
I & 0  \tag{68}\\
-\frac{1}{2} w^{T} & 1
\end{array}\right]\left[\begin{array}{cc}
\hat{Q} & 0 \\
0 & 1-\frac{1}{4} w^{T} \hat{Q} w
\end{array}\right]\left[\begin{array}{cc}
I & -\frac{1}{2} w \\
0 & 1
\end{array}\right]=J\left[\begin{array}{cc}
\hat{Q} & 0 \\
0 & c
\end{array}\right] J^{T},
$$

where $J \triangleq\left[\begin{array}{cc}I & 0 \\ -\frac{1}{2} w^{T} & 1\end{array}\right], I$ denotes the $(n-1)$-dim identity matrix and

$$
\begin{equation*}
c \triangleq 1-\frac{1}{4} w^{T} \hat{Q} w . \tag{69}
\end{equation*}
$$

It is easy to prove

$$
\begin{equation*}
\operatorname{eig}\left(A Q_{n} A^{T}\right) \subseteq(0, \infty) \tag{70}
\end{equation*}
$$

In fact, since $A$ is non-singular, we only need to prove $Q_{n} \succ 0$. According to (68), we only need to prove $\left[\begin{array}{rr}\hat{Q} & 0 \\ 0 & c\end{array}\right] \succ 0$. This follows from $\hat{Q} \succ 0$, and the fact $c=1-\frac{1}{4} w^{T} \hat{Q} w \stackrel{[67)}{>} 1-\frac{1}{3}>0$. Thus (70) is proved.
It remains to prove

$$
\begin{equation*}
\rho\left(A Q_{n} A^{T}\right)<\frac{4}{3} \tag{71}
\end{equation*}
$$

Denote $\hat{B} \triangleq \hat{A}^{T} \hat{A}$, then we can write $A^{T} A$ as

$$
A^{T} A=\left[\begin{array}{cc}
\hat{B} & w \\
w^{T} & 1
\end{array}\right] .
$$

We simplify the expression of $\rho\left(A Q_{n} A^{T}\right)=\rho\left(Q_{n} A^{T} A\right)$ as follows:

$$
\rho\left(A Q_{n} A^{T}\right) \stackrel{\boxed{68}}{=} \rho\left(J\left[\begin{array}{cc}
\hat{Q} & 0  \tag{72}\\
0 & c
\end{array}\right] J^{T} A^{T} A\right)=\rho\left(\left[\begin{array}{cc}
\hat{Q} & 0 \\
0 & c
\end{array}\right] J^{T} A^{T} A J\right) .
$$

By algebraic computation, we have
$J^{T} A^{T} A J=\left[\begin{array}{cc}I & -\frac{1}{2} w \\ 0 & 1\end{array}\right]\left[\begin{array}{cc}\hat{B} & w \\ w^{T} & 1\end{array}\right]\left[\begin{array}{cc}I & 0 \\ -\frac{1}{2} w^{T} & 1\end{array}\right]=\left[\begin{array}{cc}I & -\frac{1}{2} w \\ 0 & 1\end{array}\right]\left[\begin{array}{cc}\hat{B}-\frac{1}{2} w w^{T} & w \\ \frac{1}{2} w^{T} & 1\end{array}\right]=\left[\begin{array}{cc}\hat{B}-\frac{3}{4} w w^{T} & \frac{1}{2} w \\ \frac{1}{2} w^{T} & 1\end{array}\right]$,
thus

$$
Z \triangleq\left[\begin{array}{cc}
\hat{Q} & 0  \tag{73}\\
0 & c
\end{array}\right] J^{T} A^{T} A J=\left[\begin{array}{cc}
\hat{Q} & 0 \\
0 & c
\end{array}\right]\left[\begin{array}{cc}
\hat{B}-\frac{3}{4} w w^{T} & \frac{1}{2} w \\
\frac{1}{2} w^{T} & 1
\end{array}\right]=\left[\begin{array}{cc}
\hat{Q} \hat{B}-\frac{3}{4} \hat{Q} w w^{T} & \frac{1}{2} \hat{Q} w \\
\frac{1}{2} c w^{T} & c
\end{array}\right] .
$$

According to (72), $\operatorname{eig}\left(A Q_{n} A^{T}\right)=\operatorname{eig}(Z)$, thus to prove (71), we only need to prove $\rho(Z)<\frac{4}{3}$. Since we have proved that $\operatorname{eig}(Z)=\operatorname{eig}\left(A Q_{n} A^{T}\right) \subseteq(0, \infty)$, we only need to prove $\lambda_{\max }(Z)<4 / 3$. In the rest, we will prove that for any eigenvalue of $Z$, denoted as $\lambda$, we have

$$
\begin{equation*}
\lambda<\frac{4}{3} . \tag{74}
\end{equation*}
$$

Suppose $v \in \mathbb{R}^{n} \backslash\{0\}$ is the eigenvector corresponding to $\lambda$, i.e. $Z v=\lambda v$. Partition $v$ into $v=\left[\begin{array}{l}v_{1} \\ v_{0}\end{array}\right]$, where $v_{1} \in \mathbb{R}^{n-1}, v_{0} \in \mathbb{R}$. According to the expression of $Z$ in 73, $Z v=\lambda v$ implies

$$
\begin{align*}
\left(\hat{Q} \hat{B}-\frac{3}{4} \hat{Q} w w^{T}\right) v_{1}+\frac{1}{2} \hat{Q} w v_{0} & =\lambda v_{1},  \tag{75a}\\
\frac{1}{2} c w^{T} v_{1}+c v_{0} & =\lambda v_{0} . \tag{75b}
\end{align*}
$$

If $\lambda=c$, then (74) holds since $c=1-\frac{1}{4} \theta \leq 1$. In the following, we assume $\lambda \neq c$. An immediate consequence is

$$
v_{1} \neq 0
$$

Otherwise, assume $v_{1}=0$; then 75 implies $c v_{0}=\lambda v_{0}$, which leads to $v_{0}=0$ and thus $v=0$, a contradiction.

By (75b) we get

$$
v_{0}=\frac{c}{2(\lambda-c)} w^{T} v_{1} .
$$

Plugging into 75a, we obtain

$$
\begin{equation*}
\lambda v_{1}=\left(\hat{Q} \hat{B}-\frac{3}{4} \hat{Q} w w^{T}\right) v_{1}+\frac{1}{2} \hat{Q} w \frac{c}{2(\lambda-c)} w^{T} v_{1}=\left(\hat{Q} \hat{B}+\phi \hat{Q} w w^{T}\right) v_{1} \tag{76}
\end{equation*}
$$

where

$$
\begin{equation*}
\phi=-\frac{3}{4}+\frac{c}{4(\lambda-c)}=\frac{\lambda}{4(\lambda-c)}-1=\frac{\lambda}{4 \lambda-4+\theta}-1 . \tag{77}
\end{equation*}
$$

Here we have used the definition $c=1-\frac{1}{4} w^{T} \hat{Q} w=1-\frac{1}{4} \theta$.
Denote $\hat{\lambda} \triangleq \rho(\hat{Q} \hat{B})$, then by the assumption $\hat{\lambda}=\rho\left(\hat{Q} \hat{A}^{T} \hat{A}\right) \in(0,4 / 3)$. We prove that

$$
\lambda \leq \begin{cases}\hat{\lambda}+\phi \theta, & \phi>0  \tag{78}\\ \hat{\lambda}, & \phi \leq 0\end{cases}
$$

Since $\hat{Q} \in \mathbb{R}^{(n-1) \times(n-1)}$ is a (symmetric) positive definite matrix, there exists $U \in \mathbb{R}^{(n-1) \times(n-1)}$ such that

$$
\hat{Q}=U^{T} U
$$

Pick a positive number $g=|\phi| \theta$. By (76) we have $(g+\lambda) v_{1}=\left(\hat{Q} \hat{B}+\phi \hat{Q} w w^{T}+g I\right) v_{1}$, here $I$ denotes the identity matrix with dimension $n-1$. Consequently,

$$
\begin{equation*}
g+\lambda \in \operatorname{eig}\left(\hat{Q} \hat{B}+\phi \hat{Q} w w^{T}+g I\right)=\operatorname{eig}\left(U \hat{B} U^{T}+\phi U w w^{T} U^{T}+g I\right) \tag{79}
\end{equation*}
$$

Note that $\phi U w w^{T} U^{T}$ is a rank-one symmetric matrix with a (possibly) non-zero eigenvalue $\phi w^{T} U^{T} U w=$ $\phi w^{T} \hat{Q} w=\phi \theta$. By our definition $g=|\phi| \theta \geq \phi \theta$, which implies that $g I+\phi w^{T} U^{T} U w \succeq 0$ and

$$
\rho\left(g I+\phi w^{T} U^{T} U w\right)= \begin{cases}g+\phi \theta, & \phi>0  \tag{80}\\ g, & \phi \leq 0\end{cases}
$$

Since both $U \hat{B} U^{T}=U \hat{A}^{T} \hat{A} U^{T}$ and $\phi U w w^{T} U^{T}+g I$ are symmetric PSD (Positive Semi-Definite) matrices, (79) implies

$$
\begin{align*}
g+\lambda & \leq \rho\left(U \hat{B} U^{T}+\phi U w w^{T} U^{T}+g I\right) \\
& \leq \rho\left(U \hat{B} U^{T}\right)+\rho\left(\phi U w w^{T} U^{T}+g I\right) \\
& = \begin{cases}\hat{\lambda}+g+\phi \theta, & \phi>0, \\
\hat{\lambda}+g, & \phi \leq 0,\end{cases} \tag{81}
\end{align*}
$$

which immediately leads to 78 .
We claim that 74) follows from 78. In fact, if $\phi \leq 0$, then by 78 we have $\lambda \leq \hat{\lambda}<\frac{4}{3}$, which proves 74). Next we assume $\phi>0$, which, by the definition of $\bar{\phi}$ in (77), means

$$
\begin{equation*}
1<\frac{\lambda}{4 \lambda-4+\theta} \tag{82}
\end{equation*}
$$

If $\lambda \leq 1$, then (74) already holds; thus we can assume $\lambda>1$, which implies $4 \lambda-4>0$. Combining with the fact $\theta \geq 0$, we have

$$
1<\frac{\lambda}{4 \lambda-4+\theta}<\frac{\lambda}{4 \lambda-4},
$$

which leads to $\lambda<\frac{4}{3}$. This finishes the proof of 74.

## 8 Proof of Lemma 2 for the general case

Without loss of generality, we can assume

$$
A_{i}^{T} A_{i}=I_{d_{i} \times d_{i}}, i=1, \ldots, n
$$

To show this, let us write $M_{\sigma}, M$ as $M_{\sigma}\left(A_{1}, \ldots, A_{n}\right)$ and $M\left(A_{1}, \ldots, A_{n}\right)$ respectively, i.e. functions of the coefficient matrix $\left(A_{1}, \ldots, A_{n}\right)$. Define $\tilde{A}_{i}=A_{i}\left(A_{i}^{T} A_{i}\right)^{-\frac{1}{2}}$ and

$$
D \triangleq \operatorname{Diag}\left(\left(A_{1}^{T} A_{1}\right)^{-\frac{1}{2}}, \ldots,\left(A_{n}^{T} A_{n}\right)^{-\frac{1}{2}}, I_{N \times N}\right)
$$

It is easy to verify that $M_{\sigma}\left(A_{1}, \ldots, A_{n}\right)=D^{-1} M_{\sigma}\left(\tilde{A}_{1}, \ldots, \tilde{A}_{n}\right) D$, which implies

$$
M\left(A_{1}, \ldots, A_{n}\right)=D^{-1} M\left(\tilde{A}_{1}, \ldots, \tilde{A}_{n}\right) D .
$$

Thus $\rho\left(M\left(A_{1}, \ldots, A_{n}\right)\right)=\rho\left(M\left(\tilde{A}_{1}, \ldots, \tilde{A}_{n}\right)\right)$. In other words, normalizing $A_{i}$ to $\tilde{A}_{i}$, which satisfies $\tilde{A}_{i}^{T} \tilde{A}_{i}=I_{d_{i} \times d_{i}}$, does not change the spectral radius of $M$.

### 8.1 Proof Outline of Lemma 2 and Two Propositions

We use mathematical induction to prove Lemma2 for the $n$-block case. For the basis of the induction $(n=1)$, Lemma 2 holds since $Q A^{T} A=I_{d_{1} \times d_{1}}$. Assume Lemma 2 holds for $n-1$, we will prove Lemma 2 for $n$.
Similar to the $n$-coordinate case, we will first derive the induction formula, and then use this formula to prove the induction step.

### 8.1.1 Step 2.1: Induction formula for the $n$-block case

For any matrix $Z \in \mathbb{R}^{N \times N}$ with $n \times n$ blocks, denote $Z[i, j]$ as the $(i, j)$-th block of $Z, 1 \leq i, j \leq n$. We use the term "the $i$-th block-row" to describe the collection of blocks $Z[i, 1], \ldots, Z[i, n]$, and "the $i$-th blockcolumn" to describe the collection of blocks $Z[1, i], \ldots, Z[n, i]$. We say the row pattern of $Z$ is $\left(r_{1}, \ldots, r_{n}\right)$ and the column pattern of $Z$ is $\left(c_{1}, \ldots, c_{n}\right)$ if $Z[i, j] \in \mathbb{R}^{r_{i} \times c_{j}}, \forall 1 \leq i, j \leq n$. The multiplication of two block partitioned matrices $Z_{1}, Z_{2} \in \mathbb{R}^{N \times N}$ can be expressed using only the blocks if the column pattern of $Z_{1}$ is the same as the row pattern of $Z_{2}$.
For $k=1, \ldots, n$, we define block-permutation matrix $S_{k} \in \mathbb{R}^{N \times N}$ with $n \times n$ blocks as follows:

$$
\begin{equation*}
S_{k}[i, i] \triangleq I_{d_{i} \times d_{i}}, i=1, \ldots, k-1 ; \quad S_{k}[j, j-1] \triangleq I_{d_{j} \times d_{j}}, j=k+1, \ldots, n ; \quad S_{k}[k, n] \triangleq I_{d_{k} \times d_{k}} \tag{83}
\end{equation*}
$$

and all other entries of $S_{k}$ are set to zero. Note that the row pattern of $S_{k}$ is $\left(d_{1}, \ldots, d_{n}\right)$ and the colum pattern of $S_{k}$ is $\left(d_{1}, \ldots, d_{k-1}, d_{k+1}, \ldots, d_{n}, d_{k}\right)$. When $d_{i}=1, \forall i$, this definition reduces to the definition (48) in the $n$-coordinate case. Similar to the $n$-coordinate case, this matrix $S_{k}$ has the following block-column moving property: right multiplying a matrix with $n$ block-columns and column pattern
$\left(d_{1}, \ldots, d_{n}\right)$ by $S_{k}$ will move the $k$-th block-column to the end, resulting in a new matrix with column pattern $\left(d_{1}, \ldots, d_{k-1}, d_{k+1}, \ldots, d_{n}, d_{k}\right)$. Consequently, $S_{k}^{T}$ has the following block-row moving property: left multiplying a matrix with $n$ block-rows by $S_{k}^{T}$ will move the $k$-th block-row to the end. Note that $S_{n}$ is the identity matrix. Another property is

$$
\begin{equation*}
S_{k}^{T}=S_{k}^{-1} \tag{84}
\end{equation*}
$$

For any $k \in[n]$, define $\Gamma_{k}$ as in (44). For any $\sigma^{\prime} \in \Gamma_{k}, L_{\sigma^{\prime}} \in \mathbb{R}^{\left(N-d_{k}\right) \times\left(N-d_{k}\right)}$ is partitioned into $(n-1) \times$ $(n-1)$ blocks and the $\left(\sigma^{\prime}(i), \sigma^{(j)}\right)$-th block is defined by

$$
L_{\sigma^{\prime}}\left[\sigma^{\prime}(i), \sigma^{\prime}(j)\right] \triangleq \begin{cases}A_{\sigma^{\prime}(i)}^{T} A_{\sigma^{\prime}(j)} & i \geq j,  \tag{85}\\ 0 & i<j\end{cases}
$$

We then define $\hat{Q}_{k} \in \mathbb{R}^{\left(N-d_{k}\right) \times\left(N-d_{k}\right)}$ by

$$
\begin{equation*}
\hat{Q}_{k} \triangleq \frac{1}{\left|\Gamma_{k}\right|} \sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\sigma^{\prime}}^{-1}, k=1, \ldots, n . \tag{86}
\end{equation*}
$$

Define $W_{k}$ as the $k$-th block-column of $A^{T} A$ excluding the block $A_{k}^{T} A_{k}$, i.e.

$$
\begin{equation*}
W_{k}=\left[A_{k}^{T} A_{1}, \ldots, A_{k}^{T} A_{k-1}, A_{k}^{T} A_{k+1}, \ldots, A_{k}^{T} A_{n}\right]^{T}, \forall k \in[n] . \tag{87}
\end{equation*}
$$

With these definitions, we are ready to present the induction formula.
Proposition 3 The matrix $Q=\frac{1}{|\Gamma|} \sum_{\sigma \in \Gamma} L_{\sigma}^{-1}$, where $L_{\sigma}$ and $\Gamma$ are defined by (18) and (6) respectively, can be decomposed as follows:

$$
\begin{equation*}
Q=\frac{1}{n} \sum_{k=1}^{n} S_{k} Q_{k} S_{k}^{T} \tag{88}
\end{equation*}
$$

where

$$
Q_{k} \triangleq\left[\begin{array}{cc}
\hat{Q}_{k} & -\frac{1}{2} \hat{Q}_{k} W_{k}  \tag{89}\\
-\frac{1}{2} W_{k}^{T} \hat{Q}_{k} & I_{d_{k} \times d_{k}}
\end{array}\right],
$$

in which $\hat{Q}_{k}$ is defined by 86.
Proposition 3 is a generalization of Proposition 1 from the $n$-coordinate case to the $n$-block case, and its proof is similar to the proof of Proposition 1 (with a slight difference due to the block partition). We relegate this proof to Appendix 10

### 8.1.2 Step 2.2: Bounding eigenvalues of each $Q_{k}$

According to 88, we have

$$
A Q A^{T}=\frac{1}{n} \sum_{k=1}^{n} A S_{k} Q_{k} S_{k}^{T} A^{T}
$$

Consequently,

$$
\begin{equation*}
\frac{1}{n} \sum_{k=1}^{n} \lambda_{\min }\left(A S_{k} Q_{k} S_{k}^{T} A^{T}\right) \leq \lambda_{\min }\left(A Q A^{T}\right) \leq \lambda_{\max }\left(A Q A^{T}\right) \leq \frac{1}{n} \sum_{k=1}^{n} \lambda_{\max }\left(A S_{k} Q_{k} S_{k}^{T} A^{T}\right) \tag{90}
\end{equation*}
$$

To prove eig $\left(A Q A^{T}\right) \subseteq(0,4 / 3)$, we only need to prove for any $k=1, \ldots, n$,

$$
\begin{equation*}
\operatorname{eig}\left(A S_{k} Q_{k} S_{k}^{T} A^{T}\right)=\operatorname{eig}\left(Q_{k} S_{k}^{T} A^{T} A S_{k}\right) \subseteq(0,4 / 3) \tag{91}
\end{equation*}
$$

By the block-column moving property of $S_{k}$, we have

$$
\begin{equation*}
\bar{A}_{k} \triangleq A S_{k}=\left[\hat{A}_{k}, A_{k}\right], \tag{92}
\end{equation*}
$$

where $\hat{A}_{k} \triangleq\left[A_{1}, \ldots, A_{k-1}, A_{k+1}, \ldots, A_{n}\right]$. Note that $\hat{Q}_{k}$ only depends on the entries of $\hat{A}_{k}^{T} \hat{A}_{k} \in$ $\mathbb{R}^{\left(N-d_{k}\right) \times\left(N-d_{k}\right)}$ which has $(n-1) \times(n-1)$ blocks, thus by the induction hypothesis, we have

$$
\begin{equation*}
\operatorname{eig}\left(\hat{Q}_{k} \hat{A}_{k}^{T} \hat{A}_{k}\right) \subseteq(0,4 / 3) \tag{93}
\end{equation*}
$$

We claim that 91 follows from the induction hypothesis (93) and the expressions (92) and (89). In fact, the following proposition directly proves fil for $k=n$. If we replace $A, \hat{A}_{n}, A_{n}, \hat{Q}_{n}, Q_{n}$ by $\overline{A_{k}}, \hat{A}_{k}, A_{k}, \hat{Q}_{k}, Q_{k}$ respectively in the following proposition, we will obtain (91) for any $k$. As mentioned earlier, the desired result $\operatorname{eig}\left(A Q A^{T}\right) \subseteq(0,4 / 3)$ follows immediately from 91) and 90).

Proposition 4 Suppose $A=\left[\hat{A}_{n}, A_{n}\right] \in \mathbb{R}^{N \times N}$ is a non-singular matrix, where $\hat{A}_{n} \in \mathbb{R}^{N \times\left(N-d_{n}\right)}$, and $A_{n} \in \mathbb{R}^{N \times d_{n}}$ satisfies $A_{n}^{T} A_{n}=I_{d_{n} \times d_{n}}$. Suppose $\hat{Q}_{n} \in \mathbb{R}^{\left(N-d_{n}\right) \times\left(N-d_{n}\right)}$ is symmetric and

$$
\begin{equation*}
\operatorname{eig}\left(\hat{Q}_{n} \hat{A}_{n}^{T} \hat{A}_{n}\right) \subseteq(0,4 / 3) . \tag{94}
\end{equation*}
$$

Define

$$
W_{n} \triangleq \hat{A}_{n}^{T} A_{n} \in \mathbb{R}^{\left(N-d_{n}\right) \times d_{n}}, \quad Q_{n} \triangleq\left[\begin{array}{cc}
\hat{Q}_{n} & -\frac{1}{2} \hat{Q}_{n} W_{n}  \tag{95}\\
-\frac{1}{2} W_{n}^{T} \hat{Q}_{n} & I_{d_{n} \times d_{n}}
\end{array}\right] .
$$

Then eig $\left(A Q_{n} A^{T}\right) \subseteq\left(0, \frac{4}{3}\right)$.
Proposition 4 is a generalization of Proposition 2 from the $n$-coordinate case to the $n$-block case, and its proof is similar to the proof of Proposition 2 (with a few minor differences). The proof of Proposition 4 is given in Section 8.2

### 8.2 Proof of Proposition 4

This proof is similar to the proof of Proposition 2 for the $n$-coordinate case, with a few minor differences due to the fact $d_{n}>1$.
For simplicity, throughout this proof, we denote

$$
W \triangleq W_{n}, \hat{Q} \triangleq \hat{Q}_{n}, \hat{A} \triangleq \hat{A}_{n} .
$$

We first prove

$$
\begin{equation*}
0 \preceq \Theta \triangleq W^{T} \hat{Q} W \prec \frac{4}{3} I . \tag{96}
\end{equation*}
$$

Since eig $\left(\hat{Q} \hat{A}^{T} \hat{A}\right) \subseteq(0, \infty)$ and $\hat{A}$ is non-singular, thus $\hat{Q} \succ 0$. Then we have $\Theta=W^{T} \hat{Q} W \succeq 0$, which proves the first relation of (96). By the definition $W=\hat{A}^{T} A_{n}$ we have

$$
\begin{array}{r}
\rho(\Theta)=\rho\left(A_{n}^{T} \hat{A} \hat{Q} \hat{A}^{T} A_{n}\right)=\max _{v \in \mathbb{R}^{d_{n} \times 1,\|v\|=1}} v^{T} A_{n}^{T} \hat{A} \hat{Q} \hat{A}^{T} A_{n} v \\
\leq \rho\left(\hat{A} \hat{Q} \hat{A}^{T}\right) \max _{v \in \mathbb{R}^{d_{n} \times 1,\|v\|=1}}\left\|A_{n} v\right\|^{2}=\rho\left(\hat{A} \hat{Q} \hat{A}^{T}\right)\left\|A_{n}\right\|^{2}=\rho\left(\hat{A} \hat{Q} \hat{A}^{T}\right)<\frac{4}{3}, \tag{97}
\end{array}
$$

where the last equality is due to the assumption $A_{n}^{T} A_{n}=I$, and the last inequality is due to the assumption (94). By 97) we have $\Theta \prec \frac{4}{3} I$, thus (96) is proved.

We apply a trick that we have previously used: factorize $Q_{n}$ and change the order of multiplication. To be specific, $Q_{n}$ defined in 95 can be factorized as

$$
Q_{n}=\left[\begin{array}{cc}
I & 0  \tag{98}\\
-\frac{1}{2} W^{T} & I
\end{array}\right]\left[\begin{array}{cc}
\hat{Q} & 0 \\
0 & I-\frac{1}{4} W^{T} \hat{Q} W
\end{array}\right]\left[\begin{array}{cc}
I & -\frac{1}{2} W \\
0 & I
\end{array}\right]=J\left[\begin{array}{cc}
\hat{Q} & 0 \\
0 & C
\end{array}\right] J^{T},
$$

where $J \triangleq\left[\begin{array}{cc}I & 0 \\ -\frac{1}{2} W^{T} & I\end{array}\right], I$ in the upper left block denotes the $\left(N-d_{n}\right)$-dimensional identity matrix, $I$ in the lower right block denotes the $d_{n}$-dim identity matrix, and

$$
\begin{equation*}
C \triangleq I-\frac{1}{4} W^{T} \hat{Q} W \in \mathbb{R}^{d_{n} \times d_{n}} \tag{99}
\end{equation*}
$$

It is easy to prove

$$
\begin{equation*}
\operatorname{eig}\left(A Q_{n} A^{T}\right) \subseteq(0, \infty) \tag{100}
\end{equation*}
$$

In fact, we only need to prove $Q_{n} \succ 0$. According to (98), we only need to prove $\left[\begin{array}{cc}\hat{Q} & 0 \\ 0 & C\end{array}\right] \succ 0$. This follows

It remains to prove

$$
\begin{equation*}
\rho\left(A Q_{n} A^{T}\right)<\frac{4}{3} \tag{101}
\end{equation*}
$$

Denote $\hat{B} \triangleq \hat{A}^{T} \hat{A} \in \mathbb{R}^{\left(N-d_{n}\right) \times\left(N-d_{n}\right)}$, then we can write $A^{T} A$ as

$$
A^{T} A=\left[\begin{array}{cc}
\hat{B} & W  \tag{102}\\
W^{T} & I
\end{array}\right] .
$$

We simplify the expression of $\rho\left(A Q_{n} A^{T}\right)$ as follows:

$$
\rho\left(A Q_{n} A^{T}\right)=\rho\left(A J\left[\begin{array}{cc}
\hat{Q} & 0  \tag{103}\\
0 & C
\end{array}\right] J^{T} A^{T}\right)=\rho\left(\left[\begin{array}{cc}
\hat{Q} & 0 \\
0 & C
\end{array}\right] J^{T} A^{T} A J\right) .
$$

By algebraic computation, we have

$$
\begin{align*}
J^{T} A^{T} A J & =\left[\begin{array}{cc}
I & -\frac{1}{2} W \\
0 & I
\end{array}\right]\left[\begin{array}{cc}
\hat{B} & W \\
W^{T} & I
\end{array}\right]\left[\begin{array}{cc}
I & 0 \\
-\frac{1}{2} W^{T} & I
\end{array}\right] \\
& =\left[\begin{array}{cc}
I & -\frac{1}{2} W \\
0 & I
\end{array}\right]\left[\begin{array}{cc}
\hat{B}-\frac{1}{2} W W^{T} & W \\
\frac{1}{2} W^{T} & I
\end{array}\right]=\left[\begin{array}{cc}
\hat{B}-\frac{3}{4} W W^{T} & \frac{1}{2} W \\
\frac{1}{2} W^{T} & I
\end{array}\right], \tag{104}
\end{align*}
$$

thus

$$
Z \triangleq\left[\begin{array}{cc}
\hat{Q} & 0  \tag{105}\\
0 & C
\end{array}\right] J^{T} A^{T} A J=\left[\begin{array}{cc}
\hat{Q} & 0 \\
0 & C
\end{array}\right]\left[\begin{array}{cc}
\hat{B}-\frac{3}{4} W W^{T} & \frac{1}{2} W \\
\frac{1}{2} W^{T} & I
\end{array}\right]=\left[\begin{array}{cc}
\hat{Q} \hat{B}-\frac{3}{4} \hat{Q} W W^{T} & \frac{1}{2} \hat{Q} W \\
\frac{1}{2} C W^{T} & C
\end{array}\right] .
$$

According to 103, $\rho\left(A Q_{n} A^{T}\right)=\rho(Z)$, thus to prove we only need to prove

$$
\rho(Z)<\frac{4}{3} .
$$

Suppose $\lambda>0$ is an arbitrary eigenvalue of $Z$. In the rest, we will prove

$$
\begin{equation*}
\lambda<\frac{4}{3} . \tag{106}
\end{equation*}
$$

Suppose $v \in \mathbb{R}^{N \times 1} \backslash\{0\}$ is the eigenvector corresponding to $\lambda$, i.e. $Z v=\lambda v$. Partition $v$ into $v=\left[\begin{array}{l}v_{1} \\ v_{0}\end{array}\right]$, where $v_{1} \in \mathbb{R}^{N-d_{n}}, v_{0} \in \mathbb{R}^{d_{n}}$. According to the expression of $Z$ in 105, $Z v=\lambda v$ implies

$$
\begin{align*}
\left(\hat{Q} \hat{B}-\frac{3}{4} \hat{Q} W W^{T}\right) v_{1}+\frac{1}{2} \hat{Q} W v_{0} & =\lambda v_{1},  \tag{107a}\\
\frac{1}{2} C W^{T} v_{1}+C v_{0} & =\lambda v_{0} . \tag{107b}
\end{align*}
$$

If $\lambda I-C$ is singular, i.e. $\lambda$ is an eigenvalue of $C$, then by (96) we have $\frac{2}{3} I \prec C=1-\frac{1}{4} \Theta \preceq I$, which implies $\lambda \leq 1$, thus 106 holds. In the following, we assume

$$
\begin{equation*}
\lambda I-C \text { is non-singular. } \tag{108}
\end{equation*}
$$

An immediate consequence is

$$
v_{1} \neq 0
$$

since otherwise 107b implies $C v_{0}=\lambda v_{0}$, which combined with leads to $v_{0}=0$ and thus $v=0$, a contradiction.

By 107b we get

$$
v_{0}=\frac{1}{2}(\lambda I-C)^{-1} C W^{T} v_{1} .
$$

Plugging into 107a, we obtain

$$
\begin{equation*}
\lambda v_{1}=\left(\hat{Q} \hat{B}-\frac{3}{4} \hat{Q} W W^{T}\right) v_{1}+\frac{1}{2} \hat{Q} W \frac{1}{2}(\lambda I-C)^{-1} C W^{T} v_{1}=\left(\hat{Q} \hat{B}+\hat{Q} W \Phi W^{T}\right) v_{1} \tag{109}
\end{equation*}
$$

where

$$
\begin{align*}
\Phi \triangleq- & \frac{3}{4} I+\frac{1}{4}(\lambda I-C)^{-1} C
\end{align*}=-I+\frac{1}{4}\left[I+(\lambda I-C)^{-1} C\right] .
$$

Here we have used the definition $C=I-\frac{1}{4} W^{T} \hat{Q} W=I-\frac{1}{4} \Theta$. Since $\Theta$ is a symmetric matrix, $\Phi$ is also a symmetric matrix.

To prove 106, we consider two cases.
Case 1: $\lambda_{\max }(\Phi)>0$.
According to 110), we have

$$
\theta \in \operatorname{eig}(\Theta) \Longleftrightarrow-1+\frac{\lambda}{(4 \lambda-4)+\theta} \in \operatorname{eig}(\Phi) .
$$

By the assumption $\lambda_{\max }(\Phi)>0$ and the above relation, there exists $\theta \in \operatorname{eig}(\Theta)$ such that

$$
\begin{equation*}
-1+\frac{\lambda}{(4 \lambda-4)+\theta}>0 \tag{111}
\end{equation*}
$$

If $\lambda<1$, then (106) already holds; so we can assume $\lambda>1$. By $\Theta \succeq 0$ we have $\theta \geq 0$, thus 111 implies $1<\frac{\lambda}{(4 \lambda-4)+\theta} \leq \frac{\lambda}{4 \lambda-4}$, which leads to $\lambda<\frac{4}{3}$. Thus in Case 1 we have proved 106.

Case 2: $\lambda_{\max }(\Phi) \leq 0$, i.e. $\Phi \preceq 0$.
By the assumption (94) we have

$$
\begin{equation*}
\hat{\lambda} \triangleq \rho(\hat{Q} \hat{B})=\rho\left(\hat{Q} \hat{A}^{T} \hat{A}\right) \in(0,4 / 3) \tag{112}
\end{equation*}
$$

Since $\hat{Q} \in \mathbb{R}^{\left(N-d_{n}\right) \times\left(N-d_{n}\right)}$ is a (symmetric) positive definite matrix, there exists a non-singular matrix $U \in \mathbb{R}^{\left(N-d_{n}\right) \times\left(N-d_{n}\right)}$ such that

$$
\begin{equation*}
\hat{Q}=U^{T} U . \tag{113}
\end{equation*}
$$

Pick a positive number $g$ that is large enough (will specify how large later). By 109 we have $(g+\lambda) v_{1}=$ $\left(\hat{Q} \hat{B}+\hat{Q} W \Phi W^{T}+g I\right) v_{1}$. Consequently,

$$
\begin{align*}
g+\lambda \in \operatorname{eig}\left(\hat{Q} \hat{B}+\hat{Q} W \Phi W^{T}+g I\right) \stackrel{[113}{=} & \operatorname{eig}\left(U^{T} U \hat{B}+U^{T} U W \Phi W^{T}+g I\right)  \tag{114}\\
& =\operatorname{eig}\left(U \hat{B} U^{T}+U W \Phi W^{T} U^{T}+g I\right) .
\end{align*}
$$

Define $\Gamma \triangleq U W \Phi W^{T} U^{T} \in \mathbb{R}^{\left(N-d_{n}\right) \times\left(N-d_{n}\right)}$, then the above relation implies

$$
\begin{align*}
g+\lambda & \leq \rho\left(U \hat{B} U^{T}+\Gamma+g I\right) \\
& \leq \rho\left(U \hat{B} U^{T}\right)+\rho(\Gamma+g I)  \tag{115}\\
& =\hat{\lambda}+\rho(\Gamma+g I),
\end{align*}
$$

where the last equality is due to $\rho\left(U \hat{B} U^{T}\right)=\rho\left(\hat{A}^{T} \hat{A} U^{T} U\right)=\rho\left(\hat{A}^{T} \hat{A} \hat{Q}\right) \stackrel{[112]}{=} \hat{\lambda}$.
For any vector $v \in \mathbb{R}^{\left(N-d_{n}\right) \times\left(N-d_{n}\right)}$ we have

$$
v^{T} \Gamma v=v^{T} U W \Phi W^{T} U^{T} v=\left(W^{T} U^{T} v\right)^{T} \Phi\left(W^{T} U^{T} v\right) \leq 0
$$

where the last inequality follows from our assumption $\Phi \preceq 0$, thus $\Gamma \preceq 0$. Pick a large $g$ so that $g>\rho(\Gamma)$, then $\rho(g I+\Gamma) \leq g$. Plugging into (115), we get

$$
g+\lambda \leq \hat{\lambda}+g
$$

which implies $\lambda \leq \hat{\lambda}<\frac{4}{3}$. Thus in Case 2 we have also proved 106.
Remark: The following generalization of $\sqrt[78]{ }$ is also true:

$$
\lambda \leq \begin{cases}\hat{\lambda}+\lambda_{\max }(\Phi)\|\Theta\|, & \lambda_{\max }(\Phi)>0  \tag{116}\\ \hat{\lambda}, & \Phi \preceq 0\end{cases}
$$

The proof of $\sqrt{116}$ is a bit longer than the proof for the scalar case, and we will omit it in this paper. Note that (116) is not necessary for the proof of Proposition 4 In particular, when $\lambda_{\max }(\Phi)>0$, we do not need to use $\lambda \leq \hat{\lambda}+\lambda_{\max }(\Phi)\|\Theta\|$ to bound $\lambda$; instead, it is enough to just use $\lambda_{\max }(\Phi)>0$ to bound $\lambda$, as shown in Case 1 of the above proof.

## 9 Proof of Proposition 1, the induction formula for the $n$-coordinate case

In the $n$-coordinate case, the ambient dimension $N=n$, and the $i$-th block of $A$ is $a_{i} \in \mathbb{R}^{n \times 1}$. Denote $\left(\sigma^{\prime}, k\right)$ and $\left(k, \sigma^{\prime}\right)$ as permutations of $[n]$ that are formed by combining a permutation $\sigma^{\prime}$ and $k$. For example, if $\sigma^{\prime}=(124)$, then $\left(\sigma^{\prime}, 3\right)=(1243)$ and $\left(3, \sigma^{\prime}\right)=(3124)$.
We divide the proof into two parts. First, we present a simple formula related to the permutation matrices. Second, we apply the three levels of symmetrization.

### 9.1 Step 1. Deal with Permutation Matrices $S_{k}$

We first prove

$$
S_{k}^{T} L_{\left(\sigma^{\prime}, k\right)} S_{k}=\left[\begin{array}{cc}
L_{\sigma^{\prime}} & w_{k}  \tag{117}\\
0 & 1
\end{array}\right] .
$$

We write $L_{\sigma^{\prime}}$ as a $2 \times 2$ block matrix

$$
L_{\sigma^{\prime}}=\left[\begin{array}{ll}
Z_{11} & Z_{12}  \tag{118}\\
Z_{21} & Z_{22}
\end{array}\right]
$$

where $Z_{11} \in \mathbb{R}^{(k-1) \times(k-1)}, Z_{12} \in \mathbb{R}^{(k-1) \times(n-k)}, Z_{21} \in \mathbb{R}^{(n-k) \times(k-1)}, Z_{22} \in \mathbb{R}^{(n-k) \times(n-k)}$, and denote

$$
U_{k}=\left(a_{1}, \ldots, a_{k-1}\right) \in \mathbb{R}^{n \times(k-1)}, \quad V_{k}=\left(a_{k+1}, \ldots, a_{n}\right) \in \mathbb{R}^{n \times(n-k)}
$$

which implies

$$
\begin{align*}
w_{k} & \stackrel{\boxed{47}}{=}\left[a_{1}, \ldots, a_{k-1}, a_{k+1}, \ldots, a_{n}\right]^{T} a_{k} \\
& =\left[U_{k}, V_{k}\right]^{T} a_{k}  \tag{119}\\
& =\left[\begin{array}{c}
U_{k}^{T} a_{k} \\
V_{k}^{T} a_{k}
\end{array}\right] .
\end{align*}
$$

It is easy to verify that

$$
L_{\left(\sigma^{\prime}, k\right)}=\left[\begin{array}{ccc}
Z_{11} & U_{k}^{T} a_{k} & Z_{12} \\
0 & 1 & 0 \\
Z_{21} & V_{k}^{T} a_{k} & Z_{22}
\end{array}\right]
$$

Note that in the above expression, $\left[\begin{array}{c}U_{k}^{T} a_{k} \\ 1 \\ V_{k}^{T} a_{k}\end{array}\right]$ is the $k$ 'th column and $[0, \ldots, 0,1,0, \ldots 0] \in \mathbb{R}^{1 \times n}$ with the entry 1 in the $k$ 'th position is the $k$ 'th row. By moving the $k$ 'th column to the end and then moving the $k$ 'th row to the end, we get

$$
S_{k}^{T}\left[\begin{array}{ccc}
Z_{11} & U_{k}^{T} a_{k} & Z_{12} \\
0 & 1 & 0 \\
Z_{21} & V_{k}^{T} a_{k} & Z_{22}
\end{array}\right] S_{k}=S_{k}^{T}\left[\begin{array}{ccc}
Z_{11} & Z_{12} & U_{k}^{T} a_{k} \\
0 & 0 & 1 \\
Z_{21} & Z_{22} & V_{k}^{T} a_{k}
\end{array}\right]=\left[\begin{array}{ccc}
Z_{11} & Z_{12} & U_{k}^{T} a_{k} \\
Z_{21} & Z_{22} & V_{k}^{T} a_{k} \\
0 & 0 & 1
\end{array}\right]=\left[\begin{array}{cc}
L_{\sigma^{\prime}} & w_{k} \\
0 & 1
\end{array}\right],
$$

where the last equality follows from (118) and 119. Thus we have proved 117).

### 9.2 Step 2: Three Levels of Symmetrization

Taking the inverse of both sides of (117) and using $S_{k}^{-1}=S_{k}^{T}$, we obtain

$$
S_{k}^{T} L_{\left(\sigma^{\prime}, k\right)}^{-1} S_{k}=\left[\begin{array}{cc}
L_{\sigma^{\prime}}^{-1} & -L_{\sigma^{\prime}}^{-1} w_{k}  \tag{120}\\
0 & 1
\end{array}\right] .
$$

As the first level symmetrization, summing up (120) for all $\sigma^{\prime} \in \Gamma_{k}$ and dividing by $\left|\Gamma_{k}\right|$, we get

$$
\frac{1}{\left|\Gamma_{k}\right|} \sum_{\sigma^{\prime} \in \Gamma_{k}} S_{k}^{T} L_{\left(\sigma^{\prime}, k\right)}^{-1} S_{k}=\left[\begin{array}{cc}
\frac{1}{\left|\Gamma_{k}\right|} \sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\sigma^{\prime}}^{-1} & -\frac{1}{\left|\Gamma_{k}\right|} \sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\sigma^{\prime}}^{-1} w_{k}  \tag{121}\\
0 & 1
\end{array}\right] \stackrel{\text { 46 }}{=}\left[\begin{array}{cc}
\hat{Q}_{k} & -\hat{Q}_{k} w_{k} \\
0 & 1
\end{array}\right] .
$$

As the second level symmetrization, we can prove

$$
\frac{1}{2\left|\Gamma_{k}\right|} S_{k}^{T}\left(\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(\sigma^{\prime}, k\right)}^{-1}+\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(k, \sigma^{\prime}\right)}^{-1}\right) S_{k}=\left[\begin{array}{cc}
\hat{Q}_{k} & -\frac{1}{2} \hat{Q}_{k} w_{k}  \tag{122}\\
\frac{1}{2} w_{k}^{T} \hat{Q}_{k} & 1
\end{array}\right]=Q_{k} .
$$

In fact, by the definition of $L_{\sigma}$ in (18), it is easy to see that

$$
L_{\sigma}^{T}=L_{\bar{\sigma}},
$$

where $\bar{\sigma}$ is a "reverse permutation" of $\sigma$ that satisfies $\bar{\sigma}(i)=\sigma(n+1-i), \forall i$. Thus we have $L_{\left(\sigma^{\prime}, k\right)}^{T}=L_{\left(k, \bar{\sigma}^{\prime}\right)}^{T}$, where $\bar{\sigma}^{\prime}$ is a reverse permutation of $\sigma^{\prime}$. Summing over all $\sigma^{\prime}$, we get $\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(\sigma^{\prime}, k\right)}^{-1}=\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(k, \bar{\sigma}^{\prime}\right)}^{-T}=$ $\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(k, \sigma^{\prime}\right)}^{-T}$, where the last equality is because summing over $\overline{\sigma^{\prime}}$ is the same as summing over $\sigma^{\prime}$. Thus, we have

$$
\frac{1}{\left|\Gamma_{k}\right|} \sum_{\sigma^{\prime} \in \Gamma_{k}} S_{k}^{T} L_{\left(k, \sigma^{\prime}\right)}^{-1} S_{k}=\left(\frac{1}{\left|\Gamma_{k}\right|} \sum_{\sigma^{\prime} \in \Gamma_{k}} S_{k}^{T} L_{\left(\sigma^{\prime}, k\right)}^{-1} S_{k}\right)^{T} \stackrel{\boxed{121}}{=}\left[\begin{array}{cc}
\hat{Q}_{k} & 0 \\
-w_{k}^{T} \hat{Q}_{k} & 1
\end{array}\right] .
$$

Here we have used the fact that $\hat{Q}_{k}$ is symmetric. Combining the above relation and 121 and invoking the definition of $Q_{k}$ in (51) yields (122).
According to $\sqrt{84}$ and the fact $\left|\Gamma_{k}\right|=(n-1)$ !, we can rewrite $\sqrt{122)}$ as

$$
S_{k} Q_{k} S_{k}^{T}=\frac{1}{2(n-1)!}\left(\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(\sigma^{\prime}, k\right)}^{-1}+\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(k, \sigma^{\prime}\right)}^{-1}\right)
$$

As the third level symmetrization, summing up the above relation for $k=1, \ldots, n$ and then dividing by $n$, we get

$$
\frac{1}{n} \sum_{k=1}^{n} S_{k} Q_{k} S_{k}^{T}=\frac{1}{n} \frac{1}{2(n-1)!} \sum_{k=1}^{n}\left(\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(\sigma^{\prime}, k\right)}^{-1}+\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(k, \sigma^{\prime}\right)}^{-1}\right)=\frac{1}{2 n!} 2 \sum_{\sigma \in \Gamma} L_{\sigma}^{-1}=Q
$$

which proves 50). Q.E.D.

## 10 Proof of Proposition 3, the induction formula for the general $n$-block case

This proof is a direct extension of the proof of Proposition 1 i.e. the induction formula for the $n$-coordinate case. The major difference is Step 1 (the proof of (123), since the permutation matrix $S_{k}$ here is a blockpartitioned matrix. Step 2 is the same as the $n$-coordinate case.

### 10.1 Step 1. Deal with Permutation Matrices $S_{k}$

We will prove

$$
S_{k}^{T} L_{\left(\sigma^{\prime}, k\right)} S_{k}=\left[\begin{array}{cc}
L_{\sigma^{\prime}} & W_{k}  \tag{123}\\
0 & I
\end{array}\right]
$$

where $I$ denotes $I_{d_{k} \times d_{k}}$.
Note that $L_{\sigma^{\prime}} \in \mathbb{R}^{\left(N-d_{k}\right) \times\left(N-d_{k}\right)}$ can be viewed as a block partitioned matrix with $(n-1) \times(n-1)$ blocks, and both the row pattern and column pattern of $L_{\sigma^{\prime}}$ are $\left(d_{1}, \ldots, d_{k-1}, d_{k+1}, \ldots, d_{n}\right)$. By grouping the first $(k-1)$ block-rows and the last $(n-k)$ block-rows respectively, and grouping the first $(k-1)$ block-columns and the last $(n-k)$ block-columns respectively, $L_{\sigma^{\prime}}$ can be written as a $2 \times 2$ block matrix

$$
L_{\sigma^{\prime}}=\left[\begin{array}{ll}
Z_{11} & Z_{12}  \tag{124}\\
Z_{21} & Z_{22}
\end{array}\right]
$$

where $Z_{11} \in \mathbb{R}^{\left(d_{1}+\cdots+d_{k-1}\right) \times\left(d_{1}+\cdots+d_{k-1}\right)}, Z_{22} \in \mathbb{R}^{\left(d_{k+1}+\cdots+d_{n}\right) \times\left(d_{k+1}+\cdots+d_{n}\right)}$, and the size of $Z_{12}$ and $Z_{22}$ can be determined accordingly. We denote

$$
U_{k}=\left(A_{1}, \ldots, A_{k-1}\right) \in \mathbb{R}^{N \times\left(d_{1}+\cdots+d_{k-1}\right)}, \quad V_{k}=\left(A_{k+1}, \ldots, A_{n}\right) \in \mathbb{R}^{N \times\left(d_{k+1}+\cdots+d_{n}\right)},
$$

which implies

$$
\begin{align*}
W_{k} & \stackrel{877}{=}\left[A_{k}^{T} A_{1}, \ldots, A_{k}^{T} A_{k-1}, A_{k}^{T} A_{k+1}, \ldots, A_{k}^{T} A_{n}\right]^{T} \\
& =\left[A_{1}, \ldots, A_{k-1}, A_{k+1}, \ldots, A_{n}\right]^{T} A_{k} \\
& =\left[U_{k}, V_{k}\right]^{T} A_{k}  \tag{125}\\
& =\left[\begin{array}{c}
U_{k}^{T} A_{k} \\
V_{k}^{T} A_{k}
\end{array}\right] .
\end{align*}
$$

It is easy to verify that

$$
L_{\left(\sigma^{\prime}, k\right)}=\left[\begin{array}{ccc}
Z_{11} & U_{k}^{T} A_{k} & Z_{12} \\
0 & I_{d_{k} \times d_{k}} & 0 \\
Z_{21} & V_{k}^{T} A_{k} & Z_{22}
\end{array}\right] .
$$

Note that in the above expression, $\left[\begin{array}{c}U_{k}^{T} A_{k} \\ I_{d_{k} \times d_{k}} \\ V_{k}^{T} A_{k}\end{array}\right]$ is the $k$ 'th block-column of $L_{\left(\sigma^{\prime}, k\right)}$ and

$$
\left[0, I_{d_{k} \times d_{k}}, 0\right]=\left[0_{d_{k} \times d_{1}}, \ldots, 0_{d_{k} \times d_{k-1}}, I_{d_{k} \times d_{k}}, 0_{d_{k} \times d_{k+1}}, \ldots 0_{d_{k} \times d_{n}}\right] \in \mathbb{R}^{d_{k} \times N}
$$

with $I_{d_{k} \times d_{k}}$ being the $k^{\prime}$ th block is the $k$ 'th block-row of $L_{\left(\sigma^{\prime}, k\right)}$. By moving the $k^{\prime}$ th block-column to the end and then moving the $k$ 'th block-row to the end, we get
$S_{k}^{T}\left[\begin{array}{ccc}Z_{11} & U_{k}^{T} a_{k} & Z_{12} \\ 0 & I_{d_{k} \times d_{k}} & 0 \\ Z_{21} & V_{k}^{T} a_{k} & Z_{22}\end{array}\right] S_{k}=S_{k}^{T}\left[\begin{array}{ccc}Z_{11} & Z_{12} & U_{k}^{T} a_{k} \\ 0 & 0 & I_{d_{k} \times d_{k}} \\ Z_{21} & Z_{22} & V_{k}^{T} a_{k}\end{array}\right]=\left[\begin{array}{ccc}Z_{11} & Z_{12} & U_{k}^{T} a_{k} \\ Z_{21} & Z_{22} & V_{k}^{T} a_{k} \\ 0 & 0 & I_{d_{k} \times d_{k}}\end{array}\right]=\left[\begin{array}{cc}L_{\sigma^{\prime}} & W_{k} \\ 0 & I_{d_{k} \times d_{k}}\end{array}\right]$,
where the last equality follows from 124 and 125 . Thus we have proved 123 .

### 10.2 Step 2: Three Levels of Symmetrization

The rest of the proof of Proposition 3 is the same as the proof of Proposition 1 except a minor difference that $w_{k}$ is replaced by $W_{k}$. For completeness, we still present the proof of Step 2 in detail.
Taking the inverse of both sides of (123) and using $S_{k}^{-1}=S_{k}^{T}$, we obtain

$$
S_{k}^{T} L_{\left(\sigma^{\prime}, k\right)}^{-1} S_{k}=\left[\begin{array}{cc}
L_{\sigma^{\prime}}^{-1} & -L_{\sigma^{\prime}}^{-1} W_{k}  \tag{126}\\
0 & I
\end{array}\right] .
$$

As the first level symmetrization, summing up (126) for all $\sigma^{\prime} \in \Gamma_{k}$ and dividing by $\left|\Gamma_{k}\right|$, we get

$$
\frac{1}{\left|\Gamma_{k}\right|} \sum_{\sigma^{\prime} \in \Gamma_{k}} S_{k}^{T} L_{\left(\sigma^{\prime}, k\right)}^{-1} S_{k}=\left[\begin{array}{cc}
\frac{1}{\left|\Gamma_{k}\right|} \sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\sigma^{\prime}}^{-1} & -\frac{1}{\left|\Gamma_{k}\right|} \sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\sigma^{\prime}}^{-1} W_{k}  \tag{127}\\
0 & I
\end{array}\right] \stackrel{86}{=}\left[\begin{array}{cc}
\hat{Q}_{k} & -\hat{Q}_{k} W_{k} \\
0 & I
\end{array}\right] .
$$

As the second level symmetrization, we will prove

$$
\frac{1}{2\left|\Gamma_{k}\right|} S_{k}^{T}\left(\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(\sigma^{\prime}, k\right)}^{-1}+\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(k, \sigma^{\prime}\right)}^{-1}\right) S_{k}=\left[\begin{array}{cc}
\hat{Q}_{k} & -\frac{1}{2} \hat{Q}_{k} W_{k}  \tag{128}\\
-\frac{1}{2} W_{k}^{T} \hat{Q}_{k} & I_{d_{k} \times d_{k}}
\end{array}\right]=Q_{k}
$$

By the definition of $L_{\sigma}$ in (18), it is easy to see that

$$
L_{\sigma}^{T}=L_{\bar{\sigma}}
$$

where $\bar{\sigma}$ is a "reverse permutation" of $\sigma$ that satisfies $\bar{\sigma}(i)=\sigma(n+1-i), \forall i$. Thus we have $L_{\left(\sigma^{\prime}, k\right)}=L_{\left(k, \bar{\sigma}^{\prime}\right)}^{T}$, where $\overline{\sigma^{\prime}}$ is a reverse permutation of $\sigma^{\prime}$. Summing over all $\sigma^{\prime}$, we get $\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(\sigma^{\prime}, k\right)}^{-1}=\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(k, \bar{\sigma}^{\prime}\right)}^{-T}=$ $\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(k, \sigma^{\prime}\right)}^{-T}$, where the last equality is because summing over $\overline{\sigma^{\prime}}$ is the same as summing over $\sigma^{\prime}$. Thus, we have

$$
\frac{1}{\left|\Gamma_{k}\right|} \sum_{\sigma^{\prime} \in \Gamma_{k}} S_{k}^{T} L_{\left(k, \sigma^{\prime}\right)}^{-1} S_{k}=\left(\frac{1}{\left|\Gamma_{k}\right|} \sum_{\sigma^{\prime} \in \Gamma_{k}} S_{k}^{T} L_{\left(\sigma^{\prime}, k\right)}^{-1} S_{k}\right)^{T} \stackrel{\boxed{1127}}{=}\left[\begin{array}{cc}
\hat{Q}_{k} & 0 \\
-W_{k}^{T} \hat{Q}_{k} & 1
\end{array}\right] .
$$

Here we have used the fact that $\hat{Q}_{k}$ is symmetric. Combining the above relation and 127) and invoking the definition of $Q_{k}$ in (89) yields (128).
According to 84 and the fact $\left|\Gamma_{k}\right|=(n-1)$ !, we can rewrite 128$)$ as

$$
S_{k} Q_{k} S_{k}^{T}=\frac{1}{2(n-1)!}\left(\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(\sigma^{\prime}, k\right)}^{-1}+\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(k, \sigma^{\prime}\right)}^{-1}\right)
$$

As the third level symmetrization, summing up the above relation for $k=1, \ldots, n$ and then dividing by $n$, we get

$$
\frac{1}{n} \sum_{k=1}^{n} S_{k} Q_{k} S_{k}^{T}=\frac{1}{n} \frac{1}{2(n-1)!} \sum_{k=1}^{n}\left(\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(\sigma^{\prime}, k\right)}^{-1}+\sum_{\sigma^{\prime} \in \Gamma_{k}} L_{\left(k, \sigma^{\prime}\right)}^{-1}\right)=\frac{1}{2 n!} 2 \sum_{\sigma \in \Gamma} L_{\sigma}^{-1}=Q
$$

which proves 88. Q.E.D.


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[^1]:    ${ }^{1}$ Rigourously speaking, these two bounds are not directly comparable since the result for the randomized version only holds with high probability, while the result for the cyclic version always holds.

[^2]:    ${ }^{2}$ For the purpose of proving Theorem 2 we do not need to prove this direction. Here we present the proof since it is quite straightforward and makes the result more comprehensive.

