# Robust High-Dimensional Low-Rank Matrix Estimation: Optimal Rate and Data-Adaptive Tuning 

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#### Abstract

The matrix lasso, which minimizes a least-squared loss function with the nuclear-norm regularization, offers a generally applicable paradigm for high-dimensional low-rank matrix estimation, but its efficiency is adversely affected by heavy-tailed distributions. This paper introduces a robust procedure by incorporating a Wilcoxon-type rank-based loss function with the nuclear-norm penalty for a unified high-dimensional low-rank matrix estimation framework. It includes matrix regression, multivariate regression and matrix completion as special examples. This procedure enjoys several appealing features. First, it relaxes the distributional conditions on random errors from sub-exponential or sub-Gaussian to more general distributions and thus it is robust with substantial efficiency gain for heavy-tailed random errors. Second, as the gradient function of the rank-based loss function is completely pivotal, it overcomes the challenge of tuning parameter selection and substantially saves the computation time by using an easily simulated tuning parameter. Third, we theoretically establish non-asymptotic error bounds with a nearly-oracle rate for the new estimator. Numerical results indicate that the new estimator can be highly competitive among existing methods, especially for heavy-tailed or skewed errors.


Keywords: heavy-tailed error, high dimension, low-rank matrix, non-asymptotic bounds, robustness, tuning parameter selection

## 1. Introduction

The estimation of low-rank matrices under high-dimensional settings has received extensive attention and in-depth research in the past decade. Its applications include recommendation systems (Ramlatchan et al., 2018), image inpainting (Zheng et al., 2018), compressed sensing (Golbabaee and Vandergheynst, 2012), sensor localization (Nguyen et al., 2019) and so on. The most popular model for low-rank matrix recovery is the linear operator model

$$
\begin{equation*}
\mathbf{y}=\mathfrak{X}\left(\mathbf{X} ; \mathbf{A}_{0}\right)+\boldsymbol{\varepsilon}, \tag{1}
\end{equation*}
$$

where $\mathbf{A}_{0} \in \mathbb{R}^{m_{1} \times m_{2}}$ is the matrix of interest, which is usually assumed to have a lowdimensional intrinsic structure, $\mathbf{y} \in \mathbb{R}^{p}$ is the response, $\mathbf{X}$ is the covariate vector/matrix
which belongs to some linear space $\mathcal{L}, \boldsymbol{\varepsilon} \in \mathbb{R}^{p}$ is the random error and $\mathfrak{X}: \mathcal{L} \times \mathbb{R}^{m_{1} \times m_{2}} \rightarrow \mathbb{R}^{p}$ is a bilinear operator with respect to each argument. We assume that $\mathbf{X}$ is independent of $\varepsilon$ and $\varepsilon$ has independent elements. This model allows us to deal with several important problems in a unified manner, including matrix regression (matrix compressed sensing)(Recht et al., 2010), multivariate linear regression (Yuan et al., 2007) and matrix completion (Candès and Recht, 2009; Gross, 2011), among others. For example, when $p=1$, $\mathfrak{X}(\mathbf{X} ; \mathbf{A})=\operatorname{tr}\left(\mathbf{A}^{\top} \mathbf{X}\right)$ and $\mathbf{X} \in \mathbb{R}^{m_{1} \times m_{2}}$ is a matrix of explanatory variables, the linear operator model (1) becomes the well-studied trace regression model (Negahban and Wainwright, 2011). See Section 2.1 for a detailed discussion.

One of the most successful estimation methods is the regularization approach based on the trade-off between fitting the target matrix to the data and minimizing the model complexity, i.e., solving

$$
\begin{equation*}
\widehat{\mathbf{A}}=\arg \min _{\mathbf{A} \in \mathcal{S}}\left\{Q_{n}(\mathbf{A})+\lambda P(\mathbf{A})\right\}, \tag{2}
\end{equation*}
$$

where $\mathcal{S}$ is a convex parameter space in $\mathbb{R}^{m_{1} \times m_{2}}, Q_{n}(\mathbf{A})$ is an empirical loss function, $\lambda$ is a tuning parameter and $P(\mathbf{A})$ is an appropriate penalization function. Under this paradigm, the most popular one may be the matrix lasso, which considers a least-squared loss with the nuclear-norm penalization or its variants. The literature in this area is vast. To name a few, for the trace regression model, Negahban and Wainwright (2011) derived non-asymptotic Frobenius norm estimation bounds under the sub-Gaussian assumption on the noise. Law et al. (2021) established a nearly optimal in-sample prediction risk bound for the rank-constrained least-squares estimator under no assumptions on the design matrix. For matrix regression, Fazel et al. (2008) and Recht et al. (2010) used the matrix lasso to explore the possibility of recovering a target matrix by observing its linear projection onto chosen dictionaries. For multivariate regression, similar ideas can be found in Yuan et al. (2007), Bunea et al. (2011, 2012), Bing et al. (2019), Kong et al. (2020) and the references therein. For the problem of noisy matrix completion, Koltchinskii et al. (2011), Negahban and Wainwright (2012), Rohde and Tsybakov (2011), among others, investigated the properties of the nuclear-norm penalized least-squares. They derived estimation error bounds which are shown to match the information-theoretic lower bound up to logarithmic factors. Other related works base on paradigm (2) include Hu et al. (2020), Hu et al. (2021) and Yu et al. (2022). Another important method for low-rank matrix estimation are based on matrix factorization framework (Sun and Luo, 2016; Ma et al., 2018; Tong et al., 2021a; Zhang et al., 2022), which stand beyond the consideration of the penalized estimation framework of the current paper.

Although the quadratic loss based methods can recover the target matrix with optimal rate under sub-Gaussian random errors, it is extremely sensitive to heavy-tailed or skewed errors. Many literature adopt robust loss instead of quadratic loss to deal with heavytailed random errors, and also establish the optimal estimation rate, see She and Chen (2017), Elsener and van de Geer (2018) and Tan et al. (2022), among others. However, these papers typically focus on one model and can not handle aforementioned important low-rank recovery problem in a unified manner. Fan et al. (2021) introduced the shrinkage principle to deal with heavy-tailed random errors in several important low-rank matrix recovery problems through the trace regression model. Although their estimators achieve the same estimation error rate as Negahban and Wainwright (2011) when the random error
has bounded second moments, an additional tuning parameter, the truncation level, needs to be determined. Consequently, a cross-validation method is inevitably required to select the regularization parameter $\lambda$ and the truncation level together, which is time-consuming and lacks theoretical guarantee. More critically, the challenges of heavy-tailed error and tuning parameter selection are usually intertwined. Solutions specifically designed for only one aspect of the two challenges could leave the other aspect more unsatisfactory. To the best of our knowledge, there is no method that solves both problems simultaneously in the low-rank matrix estimation.

In this paper, we propose a robust procedure, termed as rank matrix lasso, in a unified high-dimensional low-rank matrix estimation framework. It enjoys both robustness for heavy-tailed error distributions and computational efficiency of the tuning parameter selection. Our major contributions are listed from the following three aspects. (1) From the methodology aspect, we propose a new tuning-easy robust method incorporating a Wilcoxon-type rank-based loss function with the nuclear-norm penalty for the low-rank matrix estimation model (1). (2) From the computation aspect, as the gradient function of the rank-based loss function is completely pivotal, it overcomes the challenge of tuning parameter selection and substantially saves the computation time by using an easily simulated tuning parameter. Thus, the rank matrix lasso is tuning-easy. (3) From the theory aspect, we establish non-asymptotic error bounds with a nearly-oracle rate for the new estimator in a unified high-dimensional low-rank matrix recovery framework under much weaker assumptions on covariates and random errors. We largely overcome the bounded and centralized assumption on covariates in Wang et al. (2020) which is restrictive for low-rank matrix recovery. Technical arguments for extending the existing linear regression model to our more general linear operator model are nontrivial and may be also interesting in their own rights.

The remainder of our paper is structured as follows. In Section 2, we present the proposed rank matrix lasso procedure. Section 3 studies the theoretical non-asymptotic properties of the new estimator. Numerical studies, including simulations and a real-data application, are presented in Section 4 and Section 5 respectively. Section 6 concludes the paper. Technical proofs and additional simulation results are provided in the Appendix.

Notations. Let $\mathbf{A}=\left(a_{i j}\right) \in \mathbb{R}^{m_{1} \times m_{2}}$ be a rectangular matrix. We use $\|\mathbf{A}\|_{\infty}=$ $\max _{i, j}\left|a_{i j}\right|$ to denote the $\ell_{\infty}$ norm, $\|\mathbf{A}\|_{F}$ for Frobenius norm, $\|\mathbf{A}\|_{\text {op }}$ for operator norm, and $\|\mathbf{A}\|_{1}$ for nuclear norm. For square matrix $\mathbf{A}$, denote the smallest eigenvalues of $\mathbf{A}$ by $\lambda_{\text {min }}(\mathbf{A})$. For vectors, we use $\|\cdot\|_{1}$ and $\|\cdot\|_{2}$ for the $\ell_{1}$ and $\ell_{2}$ norms, respectively. Let $\psi_{p}(x)=e^{x^{p}}-1, p \geq 1$, then the $\psi_{p}$-Orlicz norm of a random variable $X$ is defined as: $\|X\|_{\psi_{p}}=\inf \left\{t>0: \mathbb{E}\left\{\psi_{p}(|X| / t)\right\} \leq 1\right\}$. For a random vector $\mathbf{x} \in \mathbb{R}^{d}$, we define its $\psi_{p^{-}}$ Orlicz norm $\|\mathbf{x}\|_{\psi_{p}}:=\sup _{\mathbf{v} \in \mathcal{D}^{d-1}}\left\|\mathbf{v}^{\top} \mathbf{x}\right\|_{\psi_{p}}$, where $\mathcal{D}^{d-1}$ is the $d$-dimensional unit sphere. Let $\mathbf{e}_{k}(m)$ be the $k$-th $m$-dimensional unit vector.

## 2. Methodology

In this section, we first introduce the model and several examples in Section 2.1. Then, in Section 2.2, we propose a new robust estimation method. Finally, in Section 2.3, we present a data-driven approach for tuning parameter selection.

### 2.1 Model

Consider $n$ independent observations collected from the linear operator model (1)

$$
\begin{equation*}
\mathbf{y}_{i}=\mathfrak{X}\left(\mathbf{X}_{i} ; \mathbf{A}_{0}\right)+\varepsilon_{i}, \quad i=1, \ldots, n, \tag{1}
\end{equation*}
$$

where $\mathbf{y}_{i}, \mathbf{X}_{i}$ and $\varepsilon_{i}$ are the response, covariate matrix and random error for the $i$ th observation, respectively. We assume that $\mathbf{A}_{0}$ is nearly low-rank by requiring that its singular value sequence $\left\{\sigma_{i}\left(\mathbf{A}_{0}\right)\right\}_{i=1}^{m}$ decays quickly enough, where $\sigma_{i}\left(\mathbf{A}_{0}\right)$ is the $i$-th largest singular value of $\mathbf{A}_{0}$ and $m=\min \left\{m_{1}, m_{2}\right\}$. This assumption on $\mathbf{A}_{0}$ is less stringent and more natural to model the real-world problems than the exact low-rank assumption. In particular, for a parameter $q \in[0,1]$ and a positive radius $R_{q}$, we consider $\mathbf{A}_{0}$ coming from the set

$$
\mathcal{B}_{q}\left(R_{q}\right):=\left\{\mathbf{A} \in \mathbb{R}^{m_{1} \times m_{2}} \mid \sum_{i=1}^{m} \sigma_{i}^{q}(\mathbf{A}) \leq R_{q}\right\} .
$$

Note that when $q=0$, the set $\mathcal{B}_{0}\left(R_{0}\right)$ corresponds to the set of matrices with rank at most $R_{0}$. This model provides a unified high-dimensional low-rank matrix recovery framework including various cases of interest.

Example 1 (Matrix regression) The matrix regression model is a setup in which one observes random linear projections of the unknown matrix $\mathbf{A}_{0}$. Concretely speaking, we have trace inner products

$$
\begin{equation*}
y_{i}=\left\langle\mathbf{X}_{i}, \mathbf{A}_{0}\right\rangle+\varepsilon_{i}, \quad i=1, \ldots, n, \tag{2}
\end{equation*}
$$

where $\left\langle\mathbf{X}_{i}, \mathbf{A}_{0}\right\rangle=\operatorname{tr}\left(\mathbf{X}_{i}^{\top} \mathbf{A}_{0}\right), \mathbf{X}_{i} \in \mathcal{L}=\mathbb{R}^{m_{1} \times m_{2}}$ is a random matrix so that $\left\langle\mathbf{X}_{i}, \mathbf{A}_{0}\right\rangle$ is a linear projection. In the typical form of matrix regression, which is called the compressed sensing, the observation matrix $\mathbf{X}_{i}$ has independent identically distributed (i.i.d.) standard normal entries. Here we relax this restriction to general sub-Gaussian ensembles. In this case, $\mathfrak{X}\left(\mathbf{X}_{i} ; \mathbf{A}_{0}\right)=\left\langle\mathbf{X}_{i}, \mathbf{A}_{0}\right\rangle=\operatorname{tr}\left(\mathbf{X}_{i}^{\top} \mathbf{A}_{0}\right)$ and $p=1$. Moreover, model (2) includes the high-dimensional linear regression model as a special case. Let $m_{1}=m_{2}=m$, and take $\left\{\mathbf{X}_{i}\right\}_{i=1}^{n}$ and $\mathbf{A}_{0}$ to be diagonal, then $\left\langle\mathbf{X}_{i}, \mathbf{A}_{0}\right\rangle=\mathbf{x}_{i}^{\top} \boldsymbol{\theta}_{0}$, where $\mathbf{x}_{i}$ and $\boldsymbol{\theta}_{0}$ denote the vectors of diagonal elements of $\mathbf{X}_{i}$ and $\mathbf{A}_{0}$, respectively. In this special case, having a low-rank $\mathbf{A}_{0}$ is equivalent to having a sparse $\boldsymbol{\theta}_{0}$.

Example 2 (Multivariate regression) The goal of multivariate regression is to estimate a prediction function that maps covariates $\mathbf{x}_{i} \in \mathbb{R}^{m_{2}}$ to multidimensional output vectors $\mathbf{y}_{i} \in \mathbb{R}^{m_{1}}$. More specifically, consider the linear model

$$
\begin{equation*}
\mathbf{y}_{i}=\mathbf{A}_{0} \mathbf{x}_{i}+\varepsilon_{i}, \quad i=1, \ldots, n \tag{3}
\end{equation*}
$$

where $\varepsilon_{i} \in \mathbb{R}^{m_{1}}$. We can write the multivariate regression as an instance of the linear operator model (1) with $\mathbf{X}_{i}=\mathbf{x}_{i} \in \mathcal{L}=\mathbb{R}^{m_{2}}$ and $\mathfrak{X}\left(\mathbf{x}_{i} ; \mathbf{A}_{0}\right)=\mathbf{A}_{0} \mathbf{x}_{i}$.

Example 3 (Matrix completion) The matrix completion problem can be formulated into the trace inner products model (2), where the matrices $\mathbf{X}_{i} \in \mathbb{R}^{m_{1} \times m_{2}}$ are so-called masks. They are assumed to lie in

$$
\mathcal{X}=\left\{\mathbf{e}_{k}\left(m_{1}\right) \mathbf{e}_{l}^{\top}\left(m_{2}\right): 1 \leq k \leq m_{1}, 1 \leq l \leq m_{2}\right\} .
$$

We will assume that $\left\{\mathbf{X}_{i}\right\}_{i=1}^{n}$ are i.i.d. samples with some underlying distribution $\Pi$ on $\mathcal{X}$. The goal is to reconstruct all the entries of $\mathbf{A}_{0}$.

### 2.2 New Estimation Method: Rank Matrix Lasso

For the unified high-dimensional low-rank matrix recovery model (1), we consider a new estimator of $\mathbf{A}_{0}$ by minimizing the following penalized loss function

$$
\begin{equation*}
\widehat{\mathbf{A}}=\underset{\mathbf{A} \in \mathcal{S}}{\arg \min }\left\{Q_{n}(\mathbf{A})+\lambda\|\mathbf{A}\|_{1}\right\}, \tag{4}
\end{equation*}
$$

where $\|\mathbf{A}\|_{1}$ denotes the nuclear norm of $\mathbf{A}$ and $\lambda$ denotes the tuning parameter. The loss function is defined as

$$
\begin{equation*}
Q_{n}(\mathbf{A})=\frac{2(n+1)}{\sqrt{3} n(n-1)} \sum_{k=1}^{p} \sum_{i=1}^{n} \phi\left(R\left(\varepsilon_{i k}(\mathbf{A})\right)\right) \cdot \varepsilon_{i k}(\mathbf{A}), \tag{5}
\end{equation*}
$$

where $\varepsilon_{i k}(\mathbf{A})=y_{i k}-\mathbf{e}_{k}^{\top}(p) \mathfrak{X}\left(\mathbf{X}_{i} ; \mathbf{A}\right), R\left(\varepsilon_{i k}(\mathbf{A})\right)$ denotes the rank of $\varepsilon_{i k}(\mathbf{A})$ among $\varepsilon_{1 k}(\mathbf{A}), \ldots, \varepsilon_{n k}(\mathbf{A})$, and $\phi(1) \leq \phi(2) \leq \cdots \leq \phi(n)$ is a set of scores generated as $\phi(i)=$ $\varphi(i /(n+1))$ for some nondecreasing score function $\varphi(u)$ defined on the interval $(0,1)$ and standardized such that $\int_{0}^{1} \varphi(u) d u=0$ and $\int_{0}^{1} \varphi^{2}(u) d u=1$. Hereafter we denote the population version of the loss function $\mathbb{E}\left\{Q_{n}(\mathbf{A})\right\}$ by $Q(\mathbf{A})$. We will show that $\mathbf{A}_{0}$ is the minimizer of $Q(\mathbf{A})$ under some weak conditions.

The multivariate rank-based loss function $Q_{n}(\mathbf{A})$ was first proposed by Davis and McKean (1993) in low-dimensional multivariate linear model. In this paper, we consider the rank-based loss with Wilcoxon score

$$
\begin{equation*}
\varphi(u)=\sqrt{12}(u-1 / 2) \tag{6}
\end{equation*}
$$

to achieve both robustness and efficiency in the low-rank matrix recovery problem. As pointed out by She and Chen (2017), changing the squared error loss to a robust loss amounts to designing a set of multiplicative weights for $\varepsilon_{i k}(\mathbf{A})$. We can regard the rank loss as using $R\left(\varepsilon_{i k}(\mathbf{A})\right) /(n+1)-1 / 2$ to weight $\varepsilon_{i k}(\mathbf{A})$, while $\ell_{2}$ loss and $\ell_{1}$ loss use $\varepsilon_{i k}(\mathbf{A})$ and $\operatorname{sign}\left(\varepsilon_{i k}(\mathbf{A})\right)$ as weights respectively. In light of this observation, the weight $R\left(\varepsilon_{i k}(\mathbf{A})\right) /(n+1)-1 / 2$ can be regarded as a balance between weight $\varepsilon_{i k}(\mathbf{A})$ and weight $\operatorname{sign}\left(\varepsilon_{i k}(\mathbf{A})\right)$. Intuitively, outliers will not have as much impact on $R\left(\varepsilon_{i k}(\mathbf{A})\right) /(n+1)-1 / 2$ as they do for weight $\varepsilon_{i k}(\mathbf{A})$, and at the same time, information on the relative magnitude of errors can still be utilized to improve the performance of estimator comparing to the $\ell_{1}$ loss. Similar ideas also appear on the commonly used Huber loss which is a combination of $\ell_{2}$ and $\ell_{1}$ loss, but the truncation level needs to be determined.

Wang and Li (2009) considered the weighted Wilcoxon-type univariate rank-based loss with the SCAD penalty (Fan and Li, 2001) for low-dimensional linear regression. Furthermore, Wang et al. (2020) investigated the appealing features of this Wilcoxon-type univariate rank-based loss function for high-dimensional linear regression. A natural question is whether the new estimator $\widehat{\mathbf{A}}$ in (4) for high-dimensional low-rank matrix recovery problems can still inherit the merits of the rank lasso estimator for linear regressions (Wang et al., 2020). In the later sections, we name our new robust method via the penalized multivariate rank-based loss optimization with score function (6) for high-dimensional low-rank
matrix recovery problems as rank matrix lasso. We will show that the new estimator $\widehat{\mathbf{A}}$ behaves very similarly as matrix lasso for normal random errors and remains robust under heavy-tailed errors.

### 2.3 The Choice of the Tuning Parameter

It is critical to select the tuning parameter $\lambda$ for the regularization methods in a computationally efficient way as different $\lambda$ 's may produce quite different models. Traditional cross-validation (CV) or information criteria techniques are computationally inefficient to exhaustively search an appropriate value of $\lambda$. Fortunately, it was noted in Wang et al. (2020) that for high-dimensional linear regression model, the gradient function of the rank based loss function is completely pivotal (Belloni et al., 2011; Parzen et al., 1994), leading to an appealing tuning-easy property. This inspires us to consider whether the similar property can be achieved by our rank matrix lasso method. Pivotal tuning would be especially interesting in matrix cases, since it allows us to circumvent the difficulty of tuning parameter selection for high-dimensional matrix estimation problems, which are typically very time-consuming if we apply conventional selection criteria such like cross-validation.

Let $R_{i k}=\operatorname{rank}\left(\varepsilon_{i k}\right)$ be the rank of $\varepsilon_{i k}$ among $\left\{\varepsilon_{1 k}, \ldots, \varepsilon_{n k}\right\}$. Write $\xi_{i k}=2 R_{i k}-(n+1)$ for $i=1, \ldots, n, k=1, \ldots, p$. By the definition of $Q_{n}(\mathbf{A})$, direct computation yields the gradient of $Q_{n}(\mathbf{A})$ evaluated at $\mathbf{A}_{0}$,

$$
\nabla Q_{n}\left(\mathbf{A}_{0}\right)=-2\{n(n-1)\}^{-1} \sum_{k=1}^{p} \sum_{i=1}^{n} \mathbf{H}_{i k} \xi_{i k},
$$

where $\mathbf{H}_{i k} \in \mathbb{R}^{m_{1} \times m_{2}}$ and the $(a, b)$-th element of $\mathbf{H}_{i k}$ is $\mathbf{e}_{k}^{\top}(p) \mathfrak{X}\left(\mathbf{X}_{i} ; \mathbf{E}_{a b}\right)$, here $\mathbf{E}_{a b}=$ $\mathbf{e}_{a}\left(m_{1}\right) \mathbf{e}_{b}^{\top}\left(m_{2}\right) \in \mathbb{R}^{m_{1} \times m_{2}}$.

It is important to observe that $\left\{R_{1 k}, R_{2 k}, \ldots, R_{n k}\right\}$ follows the uniform distribution on the permutations of the integers $\{1,2, \ldots, n\}$. Therefore, $\nabla Q_{n}\left(\mathbf{A}_{0}\right)$ has a known distribution conditional on covariates $\mathbf{X}_{1}, \cdots, \mathbf{X}_{n}$.

By the theoretical analysis given in Section 3, conditional on the event that

$$
\begin{equation*}
\lambda \geq 2\left\|\nabla Q_{n}\left(\mathbf{A}_{0}\right)\right\|_{\mathrm{op}} \tag{7}
\end{equation*}
$$

the rank matrix lasso estimator enjoys the nearly-oracle error bound, $\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F} \propto \lambda^{1-q / 2}$. Larger $\lambda$ increases the probability of that event but will have an adverse effect on estimation accuracy. This suggests that it is desirable to choose a small $\lambda$ such that the event (7) holds with high probability. In the same spirit of Wang et al. (2020), we introduce a new variable $S_{n}=\left\|\nabla Q_{n}\left(\mathbf{A}_{0}\right)\right\|_{\mathrm{op}}$ and recommend to take $\lambda$ equal to

$$
\begin{equation*}
\lambda^{*}=2 G_{S_{n}}^{-1}\left(1-\alpha_{0}\right), \tag{8}
\end{equation*}
$$

where $G_{S_{n}}^{-1}\left(1-\alpha_{0}\right)$ denotes the $\left(1-\alpha_{0}\right)$ th quantile of the distribution of $S_{n}$ conditional on covariates $\mathbf{X}_{1}, \cdots, \mathbf{X}_{n}$. Because $S_{n}$ is distribution-free as discussed above, the $\lambda^{*}$ does not depend on the estimation of any unknown population quantity and thus can be obtained via a simulation method given $\mathbf{X}_{1}, \ldots, \mathbf{X}_{n}$.

## 3. Non-asymptotic Properties

We first present Theorem 1 that serves as a roadmap to establish the convergence rate for $\widehat{\mathbf{A}}$. Then we apply this theorem to the three specific problems, matrix regression, multivariate regression and matrix completion, in Sections 3.1-3.3, respectively, and derive explicit nonasymptotic error bounds which allow us to compare with existing works.

Our main result is given as follows. The set $\mathcal{S}$ is a convex parameter space which is determined based on the concrete settings.

Theorem 1 Suppose $\mathbf{A}_{0} \in \mathcal{B}_{q}\left(R_{q}\right) \cap \mathcal{S}$ and that the regularization parameter $\lambda$ is chosen such that $\lambda \geq 2\left\|\nabla Q_{n}\left(\mathbf{A}_{0}\right)\right\|_{\mathrm{op}}$. Suppose further that for $\lambda_{\varepsilon} \geq 0, \widetilde{\lambda} \geq 0$ and all $\mathbf{A} \in \mathcal{S}$,

$$
\left|\left\{Q_{n}(\mathbf{A})-Q(\mathbf{A})\right\}-\left\{Q_{n}\left(\mathbf{A}_{0}\right)-Q\left(\mathbf{A}_{0}\right)\right\}\right| \leq \lambda_{\varepsilon}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1}+\widetilde{\lambda} .
$$

Then for each integer $r \in\{1,2, \ldots, m\}$, the estimator $\widehat{\mathbf{A}}$ satisfies

$$
Q(\widehat{\mathbf{A}})-Q\left(\mathbf{A}_{0}\right) \leq \max \left\{12 \sqrt{2 r}\left(\lambda_{\varepsilon}+\lambda\right)\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}, 12\left(\lambda_{\varepsilon}+\lambda\right) \sum_{j=r+1}^{m} \sigma_{j}\left(\mathbf{A}_{0}\right), 3 \widetilde{\lambda}\right\} .
$$

Define the restricted set

$$
\mathcal{C}:=\left\{\boldsymbol{\Delta} \in \mathbb{R}^{m_{1} \times m_{2}} \mid\left\|\boldsymbol{\Delta}^{\prime \prime}\right\|_{1} \leq 3\left\|\boldsymbol{\Delta}^{\prime}\right\|_{1}+4 \sum_{j=r+1}^{m} \sigma_{j}\left(\mathbf{A}_{0}\right)\right\} .
$$

If $Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right) \geq \kappa\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{F}^{2}$ for some positive number $\kappa$ and all $\mathbf{A}-\mathbf{A}_{0} \in \mathcal{C}$, then we have

$$
\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F} \leq \max \left\{24 \sqrt{R_{q}}\left(\frac{\lambda+\lambda_{\varepsilon}}{\kappa}\right)^{1-q / 2}, \sqrt{\frac{3 \widetilde{\lambda}}{\kappa}}\right\}
$$

This theorem reveals that three major conditions are required to yield a convergence rate of $\widehat{\mathbf{A}}$. First, we need $\lambda$ to be greater than $2\left\|\nabla Q_{n}\left(\mathbf{A}_{0}\right)\right\|_{\mathrm{op}}$. Second, $\lambda_{\varepsilon}$ and $\widetilde{\lambda}$ are two nonrandom constants which depend on the model parameters $n, m_{1}$ and $m_{2}$. They bound the quantity $Q(\widehat{\mathbf{A}})-Q\left(\mathbf{A}_{0}\right)$ by controlling the empirical process. We need to control an empirical process to specify a proper rate of $\lambda_{\varepsilon}$ and $\widetilde{\lambda}$ by advanced empirical process techniques. At last, we need to verify that $Q(\widehat{\mathbf{A}})-Q\left(\mathbf{A}_{0}\right) \geq \kappa\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2}$ which controls the quality of minoration of the loss function by a quadratic function and strongly relates to the shape of the density function of the errors. In the restricted set $\mathcal{C}$, we take $r=$ $\#\left\{j \in\{1,2, \ldots, m\} \mid \sigma_{j}\left(\mathbf{A}_{0}\right) \geq\left(\lambda+\lambda_{\varepsilon}\right) / \kappa\right\}$ in our proof, which is known as the "effective rank" of a near low-rank matrix (Negahban and Wainwright, 2011). In specific example, it can be shown that $\lambda+\lambda_{\varepsilon}$ decay to 0 as the sample size $n$ increases, this effective rank will increase, reflecting the fact that as we obtain more samples, we can afford to estimate more of the smaller singular values of the matrix $\mathbf{A}_{0}$. The definition of $\boldsymbol{\Delta}^{\prime}$ and $\boldsymbol{\Delta}^{\prime \prime}$ are given in the Appendix A to save space.

Next we will show that, when specialized to the three examples in Section 2.1, under much weaker assumptions on random errors, the proposed robust rank matrix lasso estimators achieve the same rates as those presented in Negahban and Wainwright (2011), Negahban and Wainwright (2012), Klopp (2014). For technical convenience, we assume $\mathcal{S}=\left\{\mathbf{A} \in \mathbb{R}^{m_{1} \times m_{2}} \mid\|\mathbf{A}\|_{\infty} \leq \eta\right\}$ for some large positive constant $\eta$.

### 3.1 Matrix Regression

For the matrix regression model (2), we impose the following conditions on observation matrices and random errors.

Assumption 2 The random errors $\varepsilon_{i}$ 's are i.i.d. with density function $f(\cdot)$. Let $\zeta_{i j}=$ $\varepsilon_{i}-\varepsilon_{j}, 1 \leq i \neq j \leq n$. Denote $f^{*}(\cdot)$ as the probability density function of $\zeta_{i j}$. There exists positive constants $b_{1}$ and $b_{2}$ such that $f^{*}(0) \geq b_{1}$ and $\left|\partial f^{*}(t) / \partial t\right| \leq b_{2}$ for all $t$.

Assumption 3 The $\left\{\operatorname{vec}\left(\mathbf{X}_{i}\right)\right\}_{i=1}^{n}$ are i.i.d. sub-Gaussian vectors with $\left\|\operatorname{vec}\left(\mathbf{X}_{i}\right)\right\|_{\psi_{2}} \leq$ $\kappa_{0}<\infty$. Denote the population covariance matrix $\mathbf{J}=\operatorname{Cov}\left(\operatorname{vec}\left(\mathbf{X}_{i}\right)\right)$. Define the restricted eigenvalue

$$
\rho=\inf _{\boldsymbol{\Delta} \in \mathcal{C}, \boldsymbol{\Delta} \neq \mathbf{0}} \frac{\operatorname{vec}(\boldsymbol{\Delta})^{\top} \mathbf{J} \operatorname{vec}(\boldsymbol{\Delta})}{\operatorname{vec}(\boldsymbol{\Delta})^{\top} \operatorname{vec}(\boldsymbol{\Delta})}
$$

There exists constant $b_{3}$ such that $\rho \geq b_{3}>0$.
Assumption 4 Let $\boldsymbol{\Delta} \in \mathbb{R}^{m_{1} \times m_{2}}$, there exists a positive constant $b_{4}$ such that

$$
\frac{3 b_{1}}{2 \sqrt{2} b_{2}} \inf _{\boldsymbol{\Delta} \neq \mathbf{0}} \frac{\left(\mathbb{E}\left\langle\mathbf{X}_{1}-\mathbf{X}_{2}, \boldsymbol{\Delta}\right\rangle^{2}\right)^{3 / 2}}{\mathbb{E}\left|\left\langle\mathbf{X}_{1}-\mathbf{X}_{2}, \boldsymbol{\Delta}\right\rangle\right|^{3}} \geq b_{4}>0
$$

Remark 5 Existing works on low-rank matrix estimation usually impose sub-Gaussian distribution (Negahban and Wainwright, 2011, 2012) or bounded moment condition (Fan et al., 2021) on random errors which excludes many heavy-tailed distributions and skewed distributions such as Cauchy distribution, log-normal distribution and $\chi^{2}$ distribution. Assumption 2 relaxes such requirement to a large degree. Assumption 3 is standard to studying the error bound for matrix lasso estimators. It allows the covariance matrix $\mathbf{J}$ to be rank-degenerate when $\mathbf{A}_{0}$ is exact low rank such that $\sum_{j=r+1}^{m} \sigma_{j}\left(\mathbf{A}_{0}\right)=0$. In such case, the restricted set becomes a cone $\mathcal{C}:=\left\{\boldsymbol{\Delta} \in \mathbb{R}^{m_{1} \times m_{2}} \mid\left\|\boldsymbol{\Delta}^{\prime \prime}\right\|_{1} \leq 3\left\|\boldsymbol{\Delta}^{\prime}\right\|_{1}\right\}$. This cone completely excludes certain directions, and thus it is possible that the covariance matrix $\mathbf{J}$, while being rankdegenerate, can satisfy $\rho>0$ over the cone. A simple sufficient condition of Assumption 3 is that the smallest eigenvalue of $\mathbf{J}$ is bounded away from zero. Similar assumptions are made in recent literature on high-dimensional linear regression and low-rank matrix estimation to establish error bounds, for example, Negahban and Wainwright (2011), Fan et al. (2018), Wang et al. (2020), and Tan et al. (2022). Assumption 4 controls the quality of minoration of the loss function by a quadratic function. From our theory, we have $Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right) \geq \frac{b_{1} b_{3}}{2}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{F}^{2}$ for all $\mathbf{A}-\mathbf{A}_{0} \in \mathcal{C}$ and $\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{F} \leq b_{4}$. To guarantee that $Q(\mathbf{A})$ is sharply curved around the ground truth $\mathbf{A}_{0}$, we make the assumption $b_{4}>0$. It helps us transform the small loss difference $Q(\widehat{\mathbf{A}})-Q\left(\mathbf{A}_{0}\right)$ into small estimation error $\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}$. Assumption 4 also appears in Belloni and Chernozhukov (2011) as part of the "restricted identifiability and nonlinearity" condition. Similar conditions are used in recent literature (Gu and Zou, 2020; Zhou et al., 2023). Indeed, if the vectorized covariates $\operatorname{vec}(\mathbf{X})$ have a log-concave density, which includes many interesting distributions such as multivariate normal distributions and uniform distribution, then $b_{4} \geq 3 b_{1} /\left(2 \sqrt{2} b_{2} K\right)$ for a universal constant $K$. This follows from the fact that $\mathbb{E}|\langle\mathbf{X}, \boldsymbol{\Delta}\rangle|^{3} \leq K\left(\mathbb{E}|\langle\mathbf{X}, \boldsymbol{\Delta}\rangle|^{2}\right)^{3 / 2}$
holds for log-concave vec (X) with some universal constant $K$ by Theorem 5.22 of Lovász and Vempala (2007) and log-concavity is preserved under affine transformations and convolution; see Saumard and Wellner (2014) for a nice review of log-concavity.

It is worth noting that the above conditions can ensure that $\mathbf{A}_{0}$ is the minimizer of population version loss function $Q(\mathbf{A})$. See Lemma 22 in the Appendix. The following theorem gives the estimation error rate of $\widehat{\mathbf{A}}$.

Theorem 6 Suppose that $\mathbf{A}_{0} \in \mathcal{B}_{q}\left(R_{q}\right) \cap \mathcal{S}$ and Assumptions 2-4 hold. The regularization parameter $\lambda$ is chosen as $\lambda^{*}$. Then there exists constant $c>0$ and $C>0$ such that when $n>c\left(m_{1}+m_{2}\right) R_{q}^{2 /(2-q)}$ and $m_{1}+m_{2}>\ln \left(2 / \alpha_{0}\right)$, the estimator $\widehat{\mathbf{A}}$ satisfies

$$
\begin{equation*}
\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2} \leq C\left(\frac{\kappa_{0}}{b_{1} b_{3}}\right)^{2-q} R_{q}\left(\frac{m_{1}+m_{2}}{n}\right)^{1-q / 2} \tag{9}
\end{equation*}
$$

with probability at least $1-\alpha_{0}-2 \exp \left\{-\left(m_{1}+m_{2}\right)\right\}$.
The Frobenius norm rate here is identical to the rate established by Negahban and Wainwright (2011) under sub-Gaussian random error assumptions and also matches the minimax optimal rate of Frobenius norm established by Rohde and Tsybakov (2011). The condition $m_{1}+m_{2}>\ln \left(2 / \alpha_{0}\right)$ is very weak. For a small $\alpha_{0}=10^{-4}$, this amounts to requiring $m_{1}+m_{2}>10$. In our assumption, $b_{1}$ is kind of related to the dispersion measure of error. The smaller the dispersion of error, the larger the value that $b_{1}$ can take, resulting in smaller error bounds. This can be seen more clearly by using Gaussian error $\varepsilon \sim \mathcal{N}\left(\mu, \sigma^{2}\right)$. Now $\varepsilon_{i}-\varepsilon_{j} \sim \mathcal{N}\left(0,2 \sigma^{2}\right), b_{1}=f^{*}(0)=1 /(2 \sqrt{\pi} \sigma)$ and the corresponding estimation error bound is

$$
\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2} \leq C\left(\frac{\kappa_{0} \sigma}{b_{3}}\right)^{2-q} R_{q}\left(\frac{m_{1}+m_{2}}{n}\right)^{1-q / 2}
$$

which indicates that our results are sharp for Gaussian errors like the result in Negahban and Wainwright (2011).

Comparing our result to the result in Fan et al. (2021). Our Assumption 2 cannot be directly compared with the moment condition in Fan et al. (2021) for that the two conditions apply to different settings. Our assumption can work well for the heavy-tailed distribution without moments, like the Cauchy distribution. This is also reflected by our simulation that our method performs better than Fan et al. (2021) under the Cauchy distribution. Fan et al. (2021) do not assume the existence of error's density function and include independent but not identical distributions. In Theorem 6, as $b_{1}$ approaches infinity, our error bound approaches zero. In contrast, Fan et al. (2021) works with the moment condition $\forall i=$ $1, \ldots, n, \mathbb{E}\left|y_{i}\right|^{2 k} \leq M<\infty$ for some $k>1$, and their estimation error bound is

$$
\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2} \leq C R_{q}\left(\frac{M^{1 / k}\left(m_{1}+m_{2}\right)}{n}\right)^{1-q / 2}
$$

in which the values of $(k, M)$ quantify the effect of error dispersions.

Remark 7 In this theorem, we require that $m_{1}+m_{2}$ tends to infinity so that the optimal rate hold with high probability. When applied to $m_{1}+m_{2}$ is fixed, our theory also holds with $m_{1}+m_{2}$ replaced by $\log n$, i.e.,

$$
\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2} \leq C\left(\frac{\kappa_{0}}{b_{1} b_{3}}\right)^{2-q} R_{q}\left(\frac{\log n}{n}\right)^{1-q / 2}
$$

with probability at least $1-\alpha_{0}-2 n^{-1}$. If $m_{1}+m_{2}$ diverges faster than $n$ in practical settings, some additional assumptions about the structure of $\mathbf{A}_{0}$, such as sparsity, are necessary since the Frobenius norm rate here matches the minimax optimal rate established by Rohde and Tsybakov (2011). Recent literature Tan et al. (2022) considers the sparse reduced rank regression, they addressed this problem by adding an additional $\ell_{1}$ penalty to induce a sparse low-rank solution. Their theory implies that the $m_{1}+m_{2}$ can diverge much faster than $n$ with the sparsity condition. However, this setting is beyond the scope of our paper, and extending our approach to this case deserves future research.

Remark 8 In our theory, the number $1-\alpha_{0}$ can be regarded as the"confidence level" in the sense that our nonasymptotic bounds on the estimation error will be controlled at the optimal rate with probability close to $1-\alpha_{0}$. In Appendix $H 3$, we show the performance of our estimator under different values of $\alpha_{0}$. The estimation error is not so sensitive to the choice of $\alpha_{0}$ for different random errors and the confidence level $1-\alpha_{0} \in[0.8,0.9]$ would give good performance results in terms of balancing regularization bias with estimation variation. Our concrete recommendation for practice is to set $1-\alpha_{0}=0.8$. After the $\alpha_{0}$ is given, our estimator enjoys the minimax optimal convergence rate and our method is tuning-easy in the sense that the proposed penalization parameter is independent of the random error distribution and easier to obtain compared with other state-of-the-art methodologies.

Remark 9 In robust statistics, vertical outliers and leverage points are also worthy of attention. Our method is applicable for the vertical outliers. Consider that the vertical outliers are modeled by the Huber's $\epsilon$-contamination model (Huber, 1992), specifically, the error $\varepsilon_{i}$ follows a mixture distribution of the form $P_{\epsilon}=(1-\epsilon) P+\epsilon Q$, where $P$ is usually assumed to be some light-tailed distribution, $Q$ is an arbitrary noise distribution, and $\epsilon$ measures the strength of contamination. As long as the error $\varepsilon_{i}$ are i.i.d. from $P_{\epsilon}$, the Assumption 2 holds for a large class of $P_{\epsilon}$ that includes the distributions and their mixtures that we consider in our simulations. But our method is not applicable to situations where leverage points exist. In our method, the selection of regularization parameter $\lambda$ is related to the covariates. The existence of leverage points may cause $\lambda$ to be very large, affecting the estimation error. In the Appendix H4, we conduct simulations to compare the performance of different methods in these two cases.

For the special case, say the linear regression problem in the form of $\left\langle\mathbf{X}_{i}, \mathbf{A}_{0}\right\rangle=\mathbf{x}_{i}^{\top} \boldsymbol{\theta}_{0}$, where $\mathbf{x}_{i} \in \mathbb{R}^{m}$ is the covariate vector and $\boldsymbol{\theta}_{0} \in \mathbb{R}^{m}$ is the parameter of interest, we have the following result.

Corollary 10 Suppose that Assumptions 2-4 hold and $\sum_{i=1}^{m}\left|\theta_{0 i}\right|^{q} \leq R_{q}$, where $\boldsymbol{\theta}_{0}=$ $\left(\theta_{01}, \ldots, \theta_{0 m}\right)^{\top}$ and $0 \leq q \leq 1$. The regularization parameter $\lambda$ is chosen as $\lambda=\lambda^{*}$.

Then there exists constant $c>0$ and $C>0$ such that when $n>c R_{q}^{2 /(2-q)} \log m$ and $m>\left(3 / \alpha_{0}\right)^{1 / 3}$, the rank estimator, denoted as $\widehat{\boldsymbol{\theta}}$, satisfies

$$
\left\|\widehat{\boldsymbol{\theta}}-\boldsymbol{\theta}_{0}\right\|_{2}^{2} \leq C\left(\frac{\kappa_{0}}{b_{1} b_{3}}\right)^{2-q} R_{q}\left(\frac{\log m}{n}\right)^{1-q / 2}
$$

with probability at least $1-\alpha_{0}-3 m^{-3}$.
Due to special structures in the linear regression, the estimator $\widehat{\boldsymbol{\theta}}$ achieves a faster estimation rate than (9) in Theorem 6. This corollary is an extension of Wang et al. (2020)'s results to the sub-Gaussian design. The estimator $\widehat{\boldsymbol{\theta}}$ attains the minimax optimal rate of $\ell_{2}$ norm established by Raskutti et al. (2011).

### 3.2 Multivariate Regression

For the multivariate regression model (3), we need the following conditions.
Assumption 11 The random errors $\boldsymbol{\varepsilon}_{i}$ 's are i.i.d. with marginal density function $f_{k}(\cdot)$ for $\varepsilon_{i k}$. Let $\zeta_{i j k}=\varepsilon_{i k}-\varepsilon_{j k}, 1 \leq i \neq j \leq n, 1 \leq k \leq m_{1}$. Let $f_{k}^{*}(\cdot)$ denote the probability density function of $\zeta_{i j k}$. There exists positive constants $b_{1}$ and $b_{2}$ such that $f_{k}^{*}(0) \geq b_{1}$ and $\left|\partial f_{k}^{*}(t) / \partial t\right| \leq b_{2}$ for all $t$, uniformly in $k$.

Assumption 12 The covariates $\left\{\mathbf{x}_{i}\right\}_{i=1}^{n}$ are i.i.d. sub-Gaussian vectors with $\left\|\mathbf{x}_{i}\right\|_{\psi_{2}} \leq$ $\kappa_{0}<\infty$. Denote $\mathbf{J}=\operatorname{Cov}\left(\mathbf{x}_{i}\right)$. Define the restricted eigenvalue

$$
\rho=\inf _{\boldsymbol{\Delta} \in \mathcal{C}, \boldsymbol{\Delta} \neq \mathbf{0}} \frac{\operatorname{tr}\left(\boldsymbol{\Delta} \mathbf{J} \boldsymbol{\Delta}^{\top}\right)}{\|\boldsymbol{\Delta}\|_{F}^{2}}
$$

There exists constant $b_{3}$ such that $\rho \geq b_{3}>0$.
Assumption 13 Let $\boldsymbol{\Delta} \in \mathbb{R}^{m_{1} \times m_{2}}$, there exists a positive constant $b_{4}$ such that

$$
\frac{3 b_{1}}{2 \sqrt{2} b_{2}} \inf _{\Delta \neq \mathbf{0}} \frac{\left\{\sum_{k=1}^{m_{1}} \mathbb{E}\left|\left(\mathbf{x}_{1}-\mathbf{x}_{2}\right)^{\top} \boldsymbol{\Delta}_{k}\right|^{2}\right\}^{3 / 2}}{\sum_{k=1}^{m_{1}} \mathbb{E}\left|\left(\mathbf{x}_{1}-\mathbf{x}_{2}\right)^{\top} \boldsymbol{\Delta}_{k}\right|^{3}} \geq b_{4}>0
$$

where $\boldsymbol{\Delta}_{k}^{\top}$ is the $k$-th row of $\boldsymbol{\Delta}$.
Theorem 14 Suppose $\mathbf{A}_{0} \in \mathcal{B}_{q}\left(R_{q}\right) \cap \mathcal{S}$ and Assumptions 11-13 hold. The regularization parameter $\lambda$ is chosen such that $\lambda=\lambda^{*}$. Then there exists constant $c>0$ and $C>0$ such that when $n>c\left(m_{1}+m_{2}\right) R_{q}^{2 /(2-q)}$ and $m_{1}+m_{2}>\ln \left(2 / \alpha_{0}\right)$, the estimator $\widehat{\mathbf{A}}$ satisfies

$$
\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2} \leq C\left(\frac{\kappa_{0}}{b_{1} b_{3}}\right)^{2-q} R_{q}\left(\frac{m_{1}+m_{2}}{n}\right)^{1-q / 2}
$$

with probability at least $1-\alpha_{0}-2 \exp \left\{-\left(m_{1}+m_{2}\right)\right\}$.
Once again, this estimation rate is identical to the rate established under sub-Gaussian random error in Negahban and Wainwright (2011).

### 3.3 Matrix Completion

For the matrix completion problem (Example 3), the trace regression formulation (2) will induce an identifiability issue if we try to apply the Wilcoxon-type loss directly, due to the special structure of the masks $\mathbf{X}_{i}$. It is well-understood in classic nonparametric literature that the Wilcoxon loss is unable to extract the intercept term from a linear model (Hettmansperger and McKean, 1998). To wit, note that we have the following fact holds almost surely: $\left\langle\mathbf{X}_{i}, c \mathcal{I}\right\rangle=c$, for any $c \in \mathbb{R}$, where $\mathcal{I}=\mathbf{1 1}^{\top}$. Consequently, we can always reformulate (2) by offsetting the ground truth and introducing an intercept term

$$
y_{i}=c+\left\langle\mathbf{X}_{i}, \mathbf{A}_{0}-c \mathcal{I}\right\rangle+\varepsilon_{i} .
$$

Hence, the failure for the estimation of $c$ will cause the problem of identifiability between $\mathbf{A}_{0}$ and $\mathbf{A}_{0}-c \mathcal{I}$.

However, there is a feasible solution by introducing the Rademacher sequence $a_{i} \in$ $\{-1,+1\}$ and cope with the model

$$
\begin{equation*}
a_{i} y_{i}=\left\langle a_{i} \mathbf{X}_{i}, \mathbf{A}_{0}\right\rangle+a_{i} \varepsilon_{i}, \tag{10}
\end{equation*}
$$

where $\left\{a_{i}\right\}_{i=1}^{n}$ is independent of $\left\{\mathbf{X}_{i}, y_{i}\right\}_{i=1}^{n}$. With this simple manipulation, we can overcome the intercept issue. To see this, suppose for some constant $\mathbf{C} \in \mathbb{R}^{m_{1} \times m_{2}}$ and $c \in \mathbb{R}$, we have almost surely $\left\langle a_{i} \mathbf{X}_{i}, \mathbf{C}\right\rangle=c$. It can be easily verified that this holds only when $\mathbf{C}=\mathbf{0}$ and $c=0$, implying that no alternative offsetting formulation involving an intercept term exists for (10). Our proposed procedure is directly applicable with the pseudo observations $\left\{a_{i} y_{i}, a_{i} \mathbf{X}_{i}\right\}_{i=1}^{n}$. Thus for the matrix matrix completion problem, $Q_{n}(\mathbf{A})$ is

$$
Q_{n}(A)=\frac{2(n+1)}{\sqrt{3} n(n-1)} \sum_{i=1}^{n} \phi\left(R\left(a_{i} \varepsilon_{i}(\mathbf{A})\right)\right) \cdot a_{i} \varepsilon_{i}(\mathbf{A})
$$

where $\varepsilon_{i}(\mathbf{A})=y_{i}-\left\langle\mathbf{X}_{i}, \mathbf{A}\right\rangle$. We show that in Appendix E, (10) can guarantee that $\mathbf{A}_{0}$ is the minimizer of $Q(\mathbf{A})$ under appropriate condition.

In literature, a conventional setting for the $\mathbf{X}_{i}$ 's is that they are i.i.d sampled from the uniform distribution $\Pi$. See Rohde and Tsybakov (2011), Koltchinskii et al. (2011) and Elsener and van de Geer (2018). We consider here a more general sampling model as formulated by Klopp (2014). More precisely, let $\pi_{j k}=\mathbb{P}\left(\mathbf{X}=\mathbf{e}_{j}\left(m_{1}\right) \mathbf{e}_{k}^{\top}\left(m_{2}\right)\right)$ be the probability to observe the $(j, k)$-th entry. Denote by $C_{k}=\sum_{j=1}^{m_{1}} \pi_{j k}$ the probability to observe an element from the $k$-th column and by $R_{j}=\sum_{k=1}^{m_{2}} \pi_{j k}$ the probability to observe an element from the $j$-th row. Note that $\max _{i, j}\left(C_{i}, R_{j}\right) \geq 1 / \min \left(m_{1}, m_{2}\right)=1 / m_{2}$, where we assume $m_{1} \geq m_{2}$ without loss of generality. We impose the following assumptions on the sampling distribution and error distribution.

Assumption 15 There exists a positive constant $L \geq 1$ such that $\max _{i, j}\left(C_{i}, R_{j}\right) \leq L / m_{2}$.
Assumption 16 There exists a positive constant $\mu \geq 1$ such that $\pi_{j k} \geq\left(\mu m_{1} m_{2}\right)^{-1}$.
Then for any $\boldsymbol{\Delta} \in \mathbb{R}^{m_{1} \times m_{2}}, \mathbb{E}\left(\left\langle\boldsymbol{\Delta}, \mathbf{X}_{i}\right\rangle^{2}\right) \geq\left(\mu m_{1} m_{2}\right)^{-1}\|\boldsymbol{\Delta}\|_{F}^{2}$.

Assumption 17 The random errors $\varepsilon_{i}$ 's are i.i.d with density function $f(\cdot)$. Let $\zeta_{i j}^{-}=$ $\varepsilon_{i}-\varepsilon_{j}, \zeta_{i j}^{+}=\varepsilon_{i}+\varepsilon_{j}, 1 \leq i \neq j \leq n$. Let $f^{-}(\cdot)$ and $f^{+}(\cdot)$ denote probability density function of $\zeta_{i j}^{-}$and $\zeta_{i j}^{+}$respectively. We assume the median of $\zeta_{i j}^{+}$is 0 and there exists a positive constant $c_{1}$ such that $f^{-}(t) \geq 1 /\left(2 c_{1}^{2}\right)$ and $f^{+}(t) \geq 1 /\left(2 c_{1}^{2}\right)$ for all $|t| \leq 4 \eta$.

Write

$$
S_{n}=2\{n(n-1)\}^{-1}\left\|\sum_{i=1}^{n} a_{i} \mathbf{X}_{i} \xi_{i}\right\|_{\mathrm{op}}
$$

where $\xi_{i}=2 R_{i}-(n+1), i=1, \ldots, n$ and $\left\{R_{1}, R_{2}, \ldots, R_{n}\right\}$ follows the uniform distribution on the permutations of the integers $\{1,2, \ldots, n\}$. We recommend to take $\lambda$ equal to

$$
\lambda^{*}=2 G_{S_{n}}^{-1}\left(1-\alpha_{0}\right),
$$

where $G_{S_{n}}^{-1}\left(1-\alpha_{0}\right)$ denotes the $\left(1-\alpha_{0}\right)$-quantile of the distribution of $S_{n}$ conditional on pseudo covariates $\left\{a_{1} \mathbf{X}_{1}, \ldots, a_{n} \mathbf{X}_{n}\right\}$. We have the following theorem.

Theorem 18 Suppose $\mathbf{A}_{0} \in \mathcal{B}_{q}\left(R_{q}\right) \cap \mathcal{S}$ and Assumptions 15-17 hold. Consider the regularization parameter $\lambda=\lambda^{*}$. If $m_{1}+m_{2} \geq\left(3 / \alpha_{0}\right)^{1 / 3}$, then there exist a numerical constant $C$ such that

$$
\begin{align*}
& \left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2} \\
\leq & C \max \left\{R_{q}\left(c_{1}^{2} \mu m_{1} m_{2} \sqrt{\frac{L \log \left(m_{1}+m_{2}\right)}{n m_{2}}}\right)^{2-q}, c_{1}^{2} \eta \mu m_{1} m_{2} \sqrt{\frac{\log \left(m_{1}+m_{2}\right)}{n}}\right\} \tag{11}
\end{align*}
$$

with probability greater than $1-\beta_{0}-2\left(m_{1}+m_{2}\right)^{-2}$, where $\beta_{0}=\mathbb{P}\left(\lambda^{*} \leq 2\left\|\nabla Q_{n}\left(\mathbf{A}_{0}\right)\right\|_{\mathrm{op}}\right)$.
Remark 19 When the random errors $\varepsilon_{i}$ 's are symmetric, $a \mathbf{X}$ is independent of as (see Lemma 25 in the Appendix), which implies that $\left\|\nabla Q_{n}\left(\mathbf{A}_{0}\right)\right\|_{\text {op }}$ has the same distribution as $S_{n}$ conditional on the pseudo covariates $\left\{a_{1} \mathbf{X}_{1}, \ldots, a_{n} \mathbf{X}_{n}\right\}$. This is in accordance with previous results. However, in general cases, the distribution of $\left\|\nabla Q_{n}\left(\mathbf{A}_{0}\right)\right\|_{\text {op }}$ is unknown conditional on all the pseudo covariates $\left\{a_{1} \mathbf{X}_{1}, \ldots, a_{n} \mathbf{X}_{n}\right\}$ due to the dependence between $a \mathbf{X}$ and ac, so the pivotal tuning property is no longer valid in an exact sense. Nevertheless, we conjecture that $\beta_{0}$ in the above theorem is close to $\alpha_{0}$, say $G_{S_{n}}$ is an approximation of $\left\|\nabla Q_{n}\left(\mathbf{A}_{0}\right)\right\|_{\mathrm{op}}$. Our simulation shows that our method is still superior to other methods with such a choice of $\lambda$ in general cases.

Remark 20 When the random errors $\varepsilon_{i}$ 's are symmetric, we have $\beta_{0}=\alpha_{0}$. If $q=0, R_{0}$ becomes the rank of $\mathbf{A}_{0}$, and the error bound reduces to

$$
\frac{\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2}}{m_{1} m_{2}} \leq C \max \left\{c_{1}^{4} \mu^{2} L \frac{R_{0} m_{1} \log \left(m_{1}+m_{2}\right)}{n}, c_{1}^{2} \eta \mu \sqrt{\frac{\log \left(m_{1}+m_{2}\right)}{n}}\right\}
$$

This bound is of the same order as the one given in Theorems 7 and 10 of Klopp (2014), established under sub-Gaussian error assumptions. For the nearly low-rank case, if we
consider the matrix completion setting (i.e., $n \ll m_{1} m_{2}$ ), then the first term dominate the maximum in (11), and this rate is the same as the statistical rate of the Huber estimator and least absolute deviation estimator in Elsener and van de Geer (2018) and is only a logarithmic factor different from the minimax optimal rate established in Koltchinskii et al. (2011).

## 4. Simulation

In this section, we investigate the performance of our proposed rank matrix lasso (RML) estimator through Monte Carlo simulations. The simulation results are evaluated through 100 Monte Carlo replications. Our implementation is based on a proximal gradient algorithm which can be found in Appendix G.

Example 4.1 [Matrix regression] We firstly study the matrix regression model (2). We consider two dimensions: $m_{1}=m_{2}=m=40$ and $m_{1}=m_{2}=m=80$. The ground truth is generated by $\mathbf{A}_{0}=\mathbf{U} \mathbf{V}^{\top}$, where $\mathbf{U}$ is the first five eigenvector from the sample covariance matrix of 100 i.i.d $\mathcal{N}_{m_{1}}\left(0, \mathbf{I}_{m_{1}}\right)$ samples, $\mathbf{V}$ is the first five eigenvectors from another sample covariance of 100 i.i.d. $\mathcal{N}_{m_{2}}\left(0, \mathbf{I}_{m_{2}}\right)$ data points. The covariates are i.i.d. copies of a generic random matrix $\mathbf{X}$, which is also composed of $\mathcal{N}(0,1)$ entries(we also consider the correlated designs with varying strengths of correlation in the Appendix H5). The random errors $\varepsilon_{i}$ are sampled independently from each of the following distributions: Gaussian $\mathcal{N}(0,0.25)$, scaled Cauchy $\mathcal{C}(0,1) / 64$, and scaled and centered log-normal $\{\mathcal{L N}(0,9)-\exp (9 / 2)\} / 400$. The sample size grows from 3200 to 6400 .

We compare four nuclear norm penalization estimators: matrix lasso (Negahban and Wainwright, 2011), Robustified matrix lasso (Fan et al., 2021), regularized LAD (Elsener and van de Geer, 2018), and our rank matrix lasso. The tuning parameter of matrix lasso and regularized LAD are given by (8). In practice, (8) cannot be applied to the calculation of $\lambda^{*}$ for matrix lasso and regularized LAD for that we don't know the distribution of error. For the convenience of calculation, we assume the distribution of error is known for matrix lasso and regularized LAD. The tuning parameter $\lambda^{*}$ given by (8) is obtained by simulation based on 100 repetitions with $\alpha_{0}=0.2$. The tuning parameter of Robustified matrix lasso is given by RCV introduced in Fan et al. (2021). We use " $\ell_{2}$ ", "Robust $\ell_{2}$ ", " $\ell_{1}$ " and "RML" to denote the four methods, respectively. Note that though theoretical guarantee for the regularized LAD estimator was investigated only under the matrix completion model, we still take it as a benchmark for comparison. The logarithm values of the Frobenius norm $\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}$ for those estimators are presented in Figure 1.

All the robust estimators have much smaller statistical errors and sharper estimation results than the matrix lasso estimator under the heavy-tailed errors such as Cauchy and log-normal distribution. In particular, our RML outperforms other methods in Cauchy and log-normal cases and guarantee nearly the same performance compared with the best estimator in Gaussian case. This suggests that it is adaptive to a wide range of populations and capable of yielding a better trade-off between robustness and estimation accuracy.

In Figure 2, under the same setting as Figure 1, we compare our method with the state-of-art alternate projection method called Scaled Gradient Descent (SGD) in Tong et al. (2021a) and Scaled Subgradient Methods (SSM) in Tong et al. (2021b). The SGD and

SSM consider the $\ell_{2}$ loss and $\ell_{1}$ loss under matrix factorization framework, respectively. For the two methods, we consider three specifications on rank, $R=3,5,10$ which represent the underestimation, perfect specification and overestimation respectively. We also consider using cross validation to select rank $R$. It can be seen that the performances of SGD and SSM are sensitive to the choice of pre-specified rank $R$. Our method is either the best or is close to the best. The cross validation criterion lacks a theoretical guarantee, which is also a manifestation of the usefulness of our method.

Example 4.2 [Multivariate regression] We consider the multivariate regression model (3). Here we consider $m_{1}=m_{2}=m=40$ and $r=5$, the sample size ranges from 500 to 2000 . The ground truth $\mathbf{A}_{0}$ is generated in the same way as Example 4.1, but this example takes different covariate designs into account. As for the covariates $\mathbf{x}_{i}$, we take i.i.d. draws from a multivariate normal $\mathcal{N}_{m}(\mathbf{0}, \boldsymbol{\Sigma})$. Two choices of $\boldsymbol{\Sigma}$ are considered: (1)Identity covariance. $\boldsymbol{\Sigma}=\mathbf{I}_{m}$, which gives i.i.d. normal components for the random vectors. In this case our Assumption 12 is met with $\kappa_{0}=1$ and $b_{3}=1$. (2) Autoregressive covariance. $\boldsymbol{\Sigma}=\left(a_{i j}\right), a_{i j}=0.8^{|i-j|}$. This generates the elements of $\mathbf{x}_{i}$ from an $\operatorname{AR}(1)$ model with the coefficient fixed at 0.8. According to Grenander and Szegö (1958), in this case Assumption 12 is satisfied with $\kappa_{0}=1 / 9$ and $b_{3}=9$. We compare four methods and numerical results are presented in Figure 3. The proposed rank matrix lasso method shows a quite competitive performance within the four candidates and yields sharp accuracy in estimating the ground truth under both designs.

In addition, Figure 4 investigated the effect of heavy-tailed errors. Two settings on the matrix dimension are considered: (1) Lower dimension: $m_{1}=m_{2}=m=40$ and $n=200$; (2) Higher dimension: $m_{1}=40, m_{2}=80$, and $n=60$. The ground truth $\mathbf{A}_{0}$ is generated in the same way as Example 4.1. The covariates $\mathbf{x}_{i}$ are drawn from a $\mathcal{N}_{m_{2}}\left(\mathbf{0}, \mathbf{I}_{m_{2}}\right)$ distribution, and we simulate a sequence of independent noise vector $\varepsilon_{i} \in \mathbb{R}^{m_{1}}$, for which each component comes from a $t$ distribution with a degree of freedom $d$. Here $d$ is given by $3 k, k=1, \ldots, 6$. Figure 4 shows the results after averaging over 100 simulation runs. The performance curves of the aforementioned estimators are presented along with the degrees of freedom of the $t$ distributions. Generally as the $t$ distribution approaches the normal, better estimation accuracy can be achieved by all four estimators. When the noise has a relatively heavy tail (small $d$ ), $\ell_{2}$ loss yields a poor performance, while $\ell_{1}$ loss and the rank matrix lasso could result in a remarkable improvement. However, for the case of large $d$, the performances of $\ell_{1}$-based method would be largely compromised. In contrast, our proposed estimator still remains as competitive as $\ell_{2}$-type methods. Thus, the rank matrix lasso achieves a good balance between robustness and estimation accuracy.

Example 4.3 [Matrix completion] We study the matrix completion model in Example 3. $\mathbf{A}_{0}$ is generated similarly as before but with an additional step of normalization (divided by $\sqrt{5}$ ) such that its Forbenius norm equals 1 . The masks are i.i.d. samples from all the unit matrices, rescaled by multiplying $\sqrt{m_{1} m_{2}}$ such that the signal-to-noise ratio of the model remains a constant as the dimension grows. Again, we take $m_{1}=m_{2}=m \in\{40,80\}$. The random error takes the same form as the matrix regression example in Example 4.1, and the sample size ranges from 3200 to 6400 . The results are presented in Figure 5. The numerical results fully demonstrate the satisfactory performances of our rank matrix lasso estimators regardless of the dimension, sample size and random error, especially for log-normal random


Figure 1: Log Frobenius Errors of different estimators for matrix regression model


Figure 2: Log Frobenius Errors of different estimators for matrix regression model


Figure 3: Log Frobenius Errors of different estimators for multivariate regression model


Figure 4: Log Frobenius Errors of different estimators for multivariate regression model under $t$ noises with varying degrees of freedom
errors. The rank matrix lasso can not only be robust to heavy-tailed random errors but also perform similarly as matrix lasso under Gaussian random errors.

Moreover, the computational merit of our rank matrix lasso is that it is tuning-easy by using a simulated tuning parameter. It overcomes the challenge of tuning parameter selection and substantially saves the computation time. Next, we demonstrate the superiority of pivotal tuning using simulations. We consider the aforementioned three models, (I) matrix regression model with the same settings as Example 4.1 when $m_{1}=m_{2}=40$ and $n=3200$, (II) multivariate regression model with the same settings as Example 4.2 when $\boldsymbol{\Sigma}=\mathbf{I}_{m}, m_{1}=m_{2}=40$ and $n=500$, and (III) matrix completion with the same settings as Example 4.3 when $m_{1}=m_{2}=40$ and $n=3200$. To save the space, we only consider Fan et al. (2021)'s Robustified matrix lasso (Robust $\ell_{2}$ ) for comparison, where the robust cross-validation (CV) method is used for the tuning parameter selection. Table 1 summarizes the results including the estimation error $\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2}$, both tuning and solving computation time for different estimators. Clearly, pivotal tuning can remarkably reduce the burden of parameter tuning than the CV-based methods without scarifying estimation accuracy.

## 5. Real Data Analysis

This section is devoted to a numerical study based on the well-known Arabidopsis thaliana data, which monitors the expression levels of a group of genes contributing to the generation of isoprenoids under different experimental conditions. See Wille et al. (2004) and She and


Figure 5: Log Frobenius Errors of different estimators for matrix completion model

| Model | Error | Estimator | $\left\\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\\|_{F}^{2}$ | Rank | Tuning(s) | Solving(s) | Total(s) |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| (I) | Normal | Robust $\ell_{2}$ | 0.114(0.007) | $5(0.00)$ | 343 | 0.88 | 344 |
|  |  | RML | 0.130(0.009) | $5(0.00)$ | 0.60 | 1.37 | 1.97 |
|  | Cauchy | Robust $\ell_{2}$ | 0.105(0.007) | 12.11(2.91) | 314 | 0.89 | 315 |
|  |  | RML | <0.001 | $5(0.00)$ | 0.61 | 2.38 | 2.98 |
|  | Log-normal | Robust $\ell_{2}$ | 0.131(0.012) | 14.17(3.26) | 389 | 0.85 | 390 |
|  |  | RML | <0.001 | $5(0.00)$ | 0.63 | 3.80 | 4.43 |
| (II) | Normal | Robust $\ell_{2}$ | 0.611(0.035) | $5(0.00)$ | 3.88 | 0.03 | 3.91 |
|  |  | RML | 0.632(0.039) | $5(0.00)$ | 0.04 | 0.06 | 0.10 |
|  | Cauchy | Robust $\ell_{2}$ | 0.076(0.008) | 10.27(2.86) | 4.45 | 0.03 | 4.48 |
|  |  | RML | 0.003(0.000) | $5(0.00)$ | 0.04 | 0.08 | 0.12 |
|  | Log-normal | Robust $\ell_{2}$ | 0.323(0.074) | 21.83(4.64) | 4.88 | 0.03 | 4.91 |
|  |  | RML | $<0.001$ | $5(0.00)$ | 0.04 | 0.12 | 0.16 |
| (III) | Normal | Robust $\ell_{2}$ | 0.107(0.007) | $5(0.00)$ | 367 | 2.94 | 370 |
|  |  | RML | 0.129(0.011) | $5(0.00)$ | 0.47 | 3.19 | 3.66 |
|  | Cauchy | Robust $\ell_{2}$ | 0.103(0.007) | 13.69(2.57) | 342 | 1.83 | 344 |
|  |  | RML | 0.002(0.000) | $5(0.00)$ | 0.45 | 3.22 | 3.67 |
|  | Log-normal | Robust $\ell_{2}$ | 0.121(0.010) | 15.52(3.37) | 338 | 1.81 | 340 |
|  |  | RML | 0.048(0.006) | $5.07(0.25)$ | 0.46 | 3.16 | 3.62 |

Table 1: The comparison of estimation accuracy and computation time. (I): matrix regression model; (II): multivariate regression model; (III): matrix completion

Chen (2017) for detailed description. The study contains $n=118$ GeneChip microarray records, with the expression levels of $m_{2}=39$ genes from two upstream isoprenoid biosynthesis pathways (mevalonate and non-mevalonate) and $m_{1}=62$ genes from downstream pathways (plastoquinone, carotenoid, phytosterol and chlorophyll). We consider a multivariate regression model for the data, using genes from upstream pathways as predictors and the downstream genes as responses.

Here for the sake of comparison, we again consider the four estimators mentioned in Section 4, trained over $80 \%$ of the data, $\mathbf{Y}_{\text {train }}$, and calculate the prediction accuracy based on the remaining data serving as a test set, $\mathbf{Y}_{\text {test }}$. Concretely speaking, the accuracy for the prediction $\mathbf{Y}_{\text {pre }}$ is measured by two prediction errors, mean absolute deviation (MAD) and mean square error (MSE), as follows,

$$
\mathrm{MAD}=\frac{1}{m_{1} n_{\text {test }}}\left\|\mathbf{Y}_{\text {pre }}-\mathbf{Y}_{\text {test }}\right\|_{1,1}, \mathrm{MSE}=\frac{1}{m_{1} n_{\text {test }}}\left\|\mathbf{Y}_{\text {pre }}-\mathbf{Y}_{\text {test }}\right\|_{F}^{2}
$$

Here $\|\cdot\|_{1,1}$ simply gives the summation of the absolute values of all the entries for a given matrix. For matrix lasso, regularized LAD and our RML we apply the pivotal tuning procedure with $\alpha=0.2$ (for matrix lasso, we simply assume the error follows standard normal distribution), and for Robustified matrix lasso, we determine the tuning parameter using robust cross validation (Fan et al., 2021). We repeat the splitting step 100 times
and report the average prediction error and the standard error. The result is summarized in Table 2. We observe from Table 2 that the matrix lasso shows the highest prediction error. In contrast, the RML has consistently lower prediction errors than other methods. Compared to Robustified matrix lasso, the RML produced estimators with smaller rank and end up with a more parsimonious model, which is beneficial for some follow-up analysis such as principal component analysis or exploratory factor analysis (EFA).

| Method | $\ell_{2}$ | Robust- $\ell_{2}$ | $\ell_{1}$ | RML |
| :---: | :---: | :---: | :---: | :---: |
| MSE | $0.609(0.079)$ | $0.574(0.075)$ | $0.581(0.077)$ | $0.567(0.072)$ |
| MAD | $0.560(0.026)$ | $0.543(0.028)$ | $0.539(0.025)$ | $0.531(0.023)$ |
| Estimated rank | $3.76(0.42)$ | $8.6(0.84)$ | $3.16(0.36)$ | $4(0.20)$ |

Table 2: Prediction accuracy for the Arabidopsis thaliana generic data
Next, we perform an EFA based on the RML estimators. The analysis follows a similar manner to She and Chen (2017) which conducted a factor analysis following their additive model setup and robust reduced rank regression estimation results. Let the predicted response be $\mathbf{Y}_{\text {pre }}=\mathbf{X} \widehat{\mathbf{A}}^{\top} \in \mathbb{R}^{n \times m_{1}}$ with singular value decomposition $\widehat{\mathbf{U}} \widehat{\mathbf{D}} \widehat{\mathbf{V}}^{\top}$. Then $\widehat{\mathbf{U}}$ collects five underlying factors, and $\widehat{\mathbf{V}} \widehat{\mathbf{D}}$ records the factor loadings (coefficients) of the 62 genes in the four downstream pathways. We plot the coefficients corresponding to the first two factors in Figure 6. To identify the most significant genes for each factor, we take a same cut-off value as She and Chen (2017), which is given by $\pm m_{1}^{-1 / 2} \hat{d}_{k}$ for the $k$-th factor. Here $\hat{d}_{k}$ is the $k$-th singular value of $\widehat{\mathbf{Y}}_{\text {pre }}$, given by the $k$-th diagonal element of $\hat{\mathbf{D}}$.

Basically the factor analysis results can unveil some structural information beneath the target genes. We can see the first factor captures some joint characteristics of carotenoid and chlorophyll, and the second factor differentiates the influence of carotenoid and phytosterol, which coincide with the findings in She and Chen (2017). This reveals some group pattern in the downstream pathways that has been verified by many biological studies. For example, Trudel and Ozbun (1970) mentioned that carotenoid and chlorophyll pigments are generally interrelated, since "the two pigment systems are morphologically associated in the cell where they are attached to the same or very similar proteins in the grana or lamellae of the chloroplast", providing evidence of the group structure we summarized from the first factor loadings. It is worthy of further devotion to knit together the line of biological experiments and the insight from statistical analysis to acquire deeper understanding of these natural mechanisms.

## 6. Conclusion and Discussion

In this article, we study a linear operator model and propose a new rank matrix lasso method for high-dimensional low-rank matrix recovery which can tackle the challenges of tuning parameter selection for regularized estimator. For normal random errors, our estimator behaves very similarly as matrix lasso. It remains robust under heavy-tailed and skewed random errors in the sense that it possesses nearly optimal statistical error rates as other standard estimators under sub-Gaussian errors.


Figure 6: Factor loadings of the 62 genes from downstream pathways for the first and second factor

We conclude this article with two remarks. Firstly, it is well known that for the linear regression model, the influence function of rank-based estimators is bounded in the response space but it is unbounded in the covariate space. Hence an outlier in the covariate space can seriously impair a rank estimate. For this aspect of robustness, Fan et al. (2021) proposed several bounded moment conditions on the design matrices and showed their shrinkage principle can achieve the minimax estimation error rate. It deserves to investigate the possibility of similar extensions in our method. Secondly, while the linear operator model has already covered a wide range of problems, it's still limited in some real-life application, especially when additional structures or restrictions are imposed on the model, such as 1-bit matrix completion (Davenport et al., 2014). Hence it is attractive to explore whether this rank based method as well as the pivotal tuning property can be adapted to more general model with other side/structure information.

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## Appendix

The Appendix contains the technical proofs of all Theorems and Corollaries, the associated optimization algorithm and additional simulations results. Hereafter we let $\mathbb{E}_{\mathbb{P} \mid \mathbf{X}}$ be the conditional expectation conditioning on all observed covariates and $\mathbb{E}_{\mathbf{X}}$ be the expectation taken with respect to covariates.

## Appendix A: Proof of Theorem 1

We first give the definition of $\boldsymbol{\Delta}^{\prime}$ and $\boldsymbol{\Delta}^{\prime \prime}$, which is first proposed by Negahban and Wainwright (2011). Any matrix $\mathbf{A}_{0} \in \mathbb{R}^{m_{1} \times m_{2}}$ has a singular value decomposition of the form $\mathbf{A}_{0}=\mathbf{U D V}^{\top}$, where $\mathbf{U} \in \mathbb{R}^{m_{1} \times m_{1}}$ and $\mathbf{V} \in \mathbb{R}^{m_{2} \times m_{2}}$ are orthogonal matrices. For each integer $r \in\{1,2, \ldots, m\}$, we let $\mathbf{U}^{r} \in \mathbb{R}^{m_{1} \times r}$ and $\mathbf{V}^{r} \in \mathbb{R}^{m_{2} \times r}$ be the sub-matrices of singular vectors associated with the top $r$ singular values of $\mathbf{A}_{0}$. We then define the following two subspaces of $\mathbb{R}^{m_{1} \times m_{2}}$ :

$$
\begin{aligned}
\mathcal{A}\left(\mathbf{U}^{r}, \mathbf{V}^{r}\right): & =\left\{\boldsymbol{\Delta} \in \mathbb{R}^{m_{1} \times m_{2}} \mid \operatorname{row}(\boldsymbol{\Delta}) \subseteq \mathbf{V}^{r} \text { and } \operatorname{col}(\boldsymbol{\Delta}) \subseteq \mathbf{U}^{r}\right\} \\
\mathcal{B}\left(\mathbf{U}^{r}, \mathbf{V}^{r}\right): & \left.: \boldsymbol{\Delta} \in \mathbb{R}^{m_{1} \times m_{2}} \mid \operatorname{row}(\boldsymbol{\Delta}) \perp \mathbf{V}^{r} \text { and } \operatorname{col}(\boldsymbol{\Delta}) \perp \mathbf{U}^{r}\right\},
\end{aligned}
$$

where row $(\boldsymbol{\Delta}) \subseteq \mathbb{R}^{m_{2}}$ and $\operatorname{col}(\boldsymbol{\Delta}) \subseteq \mathbb{R}^{m_{1}}$ denote the row space and column space, respectively, of the matrix $\boldsymbol{\Delta}$. Let $\Pi_{\mathcal{B}}$ denote the projection operator onto the subspace $\mathcal{B}$, and define $\boldsymbol{\Delta}^{\prime \prime}=\Pi_{\mathcal{B}}(\boldsymbol{\Delta})$ and $\boldsymbol{\Delta}^{\prime}=\boldsymbol{\Delta}-\boldsymbol{\Delta}^{\prime \prime}$.

Before proving the main theorem, we state a useful lemma.
Lemma 21 Let $\mathbf{U}^{r} \in \mathbb{R}^{m_{1} \times r}$ and $\mathbf{V}^{r} \in \mathbb{R}^{m_{2} \times r}$ be matrices consisting of the top $r$ left and right singular vectors of $\mathbf{A}_{0}$, respectively. Then there exists a matrix decomposition $\boldsymbol{\Delta}=\boldsymbol{\Delta}^{\prime}+\boldsymbol{\Delta}^{\prime \prime}$ of the error $\boldsymbol{\Delta}=\widehat{\mathbf{A}}-\mathbf{A}_{0}$ such that:
(a) the matrix $\boldsymbol{\Delta}^{\prime}$ satisfies $\left\|\boldsymbol{\Delta}^{\prime}\right\|_{F} \leq\|\boldsymbol{\Delta}\|_{F}$ and the constraint $\operatorname{rank}\left(\boldsymbol{\Delta}^{\prime}\right) \leq 2 r$;
(b) if $\lambda \geq 2\left\|\nabla Q_{n}\left(\mathbf{A}_{0}\right)\right\|_{\text {op }}$, then the nuclear norm of $\boldsymbol{\Delta}^{\prime \prime}$ is bounded as $\left\|\boldsymbol{\Delta}^{\prime \prime}\right\|_{1} \leq$ $3\left\|\boldsymbol{\Delta}^{\prime}\right\|_{1}+4 \sum_{j=r+1}^{m} \sigma_{j}\left(\mathbf{A}_{0}\right)$.

The proof of this lemma is the same as the proof of Lemma 1 in Negahban and Wainwright (2011), thus we omit it.
Proof [Proof of Theorem 1] Using the fact that $\widehat{\mathbf{A}}$ is the minimizer of the objective function, we have

$$
Q_{n}(\widehat{\mathbf{A}})+\lambda\|\widehat{\mathbf{A}}\|_{1} \leq Q_{n}\left(\mathbf{A}_{0}\right)+\lambda\left\|\mathbf{A}_{0}\right\|_{1}
$$

From this and the assumption, we obtain the following inequality

$$
\begin{aligned}
Q(\widehat{\mathbf{A}})-Q\left(\mathbf{A}_{0}\right) & \leq-\left[\left\{Q_{n}(\widehat{\mathbf{A}})-Q(\widehat{\mathbf{A}})\right\}-\left\{Q_{n}\left(\mathbf{A}_{0}\right)-Q\left(\mathbf{A}_{0}\right)\right\}\right]+\lambda\left\|\mathbf{A}_{0}\right\|_{1}-\lambda\|\widehat{\mathbf{A}}\|_{1} \\
& \leq \lambda_{\varepsilon}\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{1}+\widetilde{\lambda}+\lambda\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{1} \\
& =\left(\lambda_{\varepsilon}+\lambda\right)\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{1}+\widetilde{\lambda} .
\end{aligned}
$$

Combining with Lemma 21, we have that

$$
\begin{aligned}
Q(\widehat{\mathbf{A}})-Q\left(\mathbf{A}_{0}\right) & \leq\left(\lambda_{\varepsilon}+\lambda\right)\left\{4 \sqrt{2 r}\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}+4 \sum_{j=r+1}^{m} \sigma_{j}\left(\mathbf{A}_{0}\right)\right\}+\tilde{\lambda} \\
& \leq \max \left\{12 \sqrt{2 r}\left(\lambda_{\varepsilon}+\lambda\right)\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}, 12\left(\lambda_{\varepsilon}+\lambda\right) \sum_{j=r+1}^{m} \sigma_{j}\left(\mathbf{A}_{0}\right), 3 \widetilde{\lambda}\right\}
\end{aligned}
$$

By Lemma 21, we have $\widehat{\mathbf{A}}-\mathbf{A}_{0} \in \mathcal{C}$, which implies that $Q(\widehat{\mathbf{A}})-Q\left(\mathbf{A}_{0}\right) \geq \kappa\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2}$. Then it follows that

$$
\kappa\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2} \leq \max \left\{12 \sqrt{2 r}\left(\lambda_{\varepsilon}+\lambda\right)\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}, 12\left(\lambda_{\varepsilon}+\lambda\right) \sum_{j=r+1}^{m} \sigma_{j}\left(\mathbf{A}_{0}\right), 3 \widetilde{\lambda}\right\}
$$

which further implies

$$
\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F} \leq \max \left\{\frac{24\left(\lambda+\lambda_{\varepsilon}\right) \sqrt{r}}{\kappa},\left\{\frac{12\left(\lambda+\lambda_{\varepsilon}\right) \sum_{j=r+1}^{m} \sigma_{j}\left(\mathbf{A}_{0}\right)}{\kappa}\right\}^{1 / 2}, \sqrt{\frac{3 \widetilde{\lambda}}{\kappa}}\right\} .
$$

For a threshold $\tau>0$, we choose $r=\#\left\{j \in\{1,2, \ldots, m\} \mid \sigma_{j}\left(\mathbf{A}_{0}\right) \geq \tau\right\}$. Then it follows that

$$
\sum_{j=r+1}^{m} \sigma_{j}\left(\mathbf{A}_{0}\right) \leq \tau \sum_{j=r+1}^{m} \frac{\sigma_{j}\left(\mathbf{A}_{0}\right)}{\tau} \leq \tau \sum_{j=r+1}^{m}\left(\frac{\sigma_{j}\left(\mathbf{A}_{0}\right)}{\tau}\right)^{q} \leq \tau^{1-q} \sum_{j=r+1}^{m} \sigma_{j}\left(\mathbf{A}_{0}\right)^{q} \leq \tau^{1-q} R_{q}
$$

On the other hand, $R_{q} \geq \sum_{j=1}^{r} \sigma_{j}\left(\mathbf{A}_{0}\right)^{q} \geq r \tau^{q}$, so $r \leq R_{q} \tau^{-q}$. Choose $\tau=\left(\lambda+\lambda_{\varepsilon}\right) / \kappa$ yields that

$$
\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F} \leq \max \left\{24 \sqrt{R_{q}}\left(\frac{\lambda+\lambda_{\varepsilon}}{\kappa}\right)^{1-q / 2}, \sqrt{\frac{3 \widetilde{\lambda}}{\kappa}}\right\}
$$

Before proceeding to the proof of Theorems 6,14 , and 18 , we give an equivalent form for rank-based loss with Wilcoxon score $Q_{n}(\mathbf{A})$, which is useful in our proof. Denote $\sum_{i \neq j}:=\sum_{i=1}^{n} \sum_{j=1, j \neq i}^{n}$, we have

$$
\begin{aligned}
Q_{n}(\mathbf{A}) & =\frac{2(n+1)}{\sqrt{3} n(n-1)} \sum_{k=1}^{p} \sum_{i=1}^{n} \phi\left(R\left(\varepsilon_{i k}(\mathbf{A})\right)\right) \cdot \varepsilon_{i k}(\mathbf{A}) \\
& =\{n(n-1)\}^{-1} \sum_{k=1}^{p} \sum_{i \neq j}\left|\varepsilon_{i k}(\mathbf{A})-\varepsilon_{j k}(\mathbf{A})\right| \\
& =\{n(n-1)\}^{-1} \sum_{k=1}^{p} \sum_{i \neq j}\left|\left\{y_{i k}-\mathbf{e}_{k}^{\top}(p) \mathfrak{X}\left(\mathbf{X}_{i} ; \mathbf{A}\right)\right\}-\left\{y_{j k}-\mathbf{e}_{k}^{\top}(p) \mathfrak{X}\left(\mathbf{X}_{j} ; \mathbf{A}\right)\right\}\right|,
\end{aligned}
$$

where $\phi(i)=\varphi(i /(n+1))$ and $\varphi(u)=\sqrt{12}(u-1 / 2)$.

## Appendix B: Proof of Theorem 6

Combining the following lemma with Theorem 1, we can complete the proof of Theorem 6.
Lemma 22 Suppose that Assumptions 2-4 hold. Suppose $\mathbf{A}_{0} \in \mathcal{B}_{q}\left(R_{q}\right) \cap \mathcal{S}$. Then we have the following conclusions.
(a) Define for all $M>0$

$$
\mathcal{Z}_{M}:=\sup _{\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1} \leq M}\left|\left\{Q_{n}(\mathbf{A})-Q(\mathbf{A})\right\}-\left\{Q_{n}\left(\mathbf{A}_{0}\right)-Q\left(\mathbf{A}_{0}\right)\right\}\right| .
$$

Then there exists a universal constant $C_{1}>0$, with $\lambda_{\varepsilon}^{\prime}=8 C_{1} \kappa_{0} \sqrt{\frac{m_{1}+m_{2}}{n}}$, such that

$$
\mathcal{Z}_{M} \leq \lambda_{\varepsilon}^{\prime} M
$$

with probability at least $1-2 \exp \left(-4\left(m_{1}+m_{2}\right)\right)$. Moreover, let $\lambda_{\varepsilon}=2 \lambda_{\varepsilon}^{\prime}$ and $\widetilde{\lambda}=$ $2 \frac{m_{1}+m_{2}}{n} \lambda_{\varepsilon}^{\prime}$, we have, for any $\mathbf{A} \in \mathcal{S}$

$$
\left|\left\{Q_{n}(\mathbf{A})-Q(\mathbf{A})\right\}-\left\{Q_{n}\left(\mathbf{A}_{0}\right)-Q\left(\mathbf{A}_{0}\right)\right\}\right| \leq \lambda_{\varepsilon}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1}+\widetilde{\lambda}
$$

with probability at least $1-2 \exp \left(-2\left(m_{1}+m_{2}\right)\right)$.
(b) $\lambda^{*}>2\left\|\nabla Q_{n}\left(\mathbf{A}_{0}\right)\right\|_{\text {op }}$ with probability $1-\alpha_{0}$. If $m_{1}+m_{2}>\ln \left(2 / \alpha_{0}\right)$, then there exists constant $C_{2}>0$ such that $\lambda^{*} \leq 4 C_{2} \kappa_{0} \sqrt{\frac{m_{1}+m_{2}}{n}}$ with probability at least $1-$ $\exp \left(-\left(m_{1}+m_{2}\right)\right)$.
(c) $\mathbf{A}_{0}$ is the minimizer of $Q(\mathbf{A})$. Take $\kappa=b_{1} b_{3} / 2$, if $\frac{16 \sqrt{2} b_{4}+16}{b_{1} b_{3} b_{4}^{2}}\left(\frac{\lambda+\lambda_{\varepsilon}}{\kappa}\right)^{1-q / 2} R_{q}^{1 / 2}<1$, then we have

$$
Q(\widehat{\mathbf{A}})-Q\left(\mathbf{A}_{0}\right) \geq \frac{b_{1} b_{3}}{2}\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2}
$$

Remark 23 Part (a) and (b) implies that $\left(\lambda^{*}+\lambda_{\varepsilon}\right)^{1-q / 2} R_{q}^{1 / 2} \lesssim \kappa_{0}\left(\frac{m_{1}+m_{2}}{n}\right)^{1 / 2-q / 4} R_{q}^{1 / 2}$. A sufficient condition of the assumption in (c) is $n \gtrsim\left(m_{1}+m_{2}\right) R_{q}^{2 /(2-q)}$, where the symbol $\gtrsim$ means that the inequality holds up to a multiplicative numerical constant.

Proof [Proof of Lemma 22] (a) The proof of the first conclusion is based on the Markov inequality, on the symmetrization theorem (Van Der Vaart and Wellner, 1996), on the contraction theorem (Ledoux and Talagrand, 2013) and on the dual norm inequality. Write $h\left(\varepsilon_{i}, \varepsilon_{j}\right)=\left|\left(\varepsilon_{i}-\varepsilon_{j}\right)-\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle\right|-\left|\varepsilon_{i}-\varepsilon_{j}\right|$,

$$
\mathcal{Z}_{M}=\sup _{\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1} \leq M}\left|\frac{1}{n(n-1)} \sum_{i \neq j}\left\{h\left(\varepsilon_{i}, \varepsilon_{j}\right)-\mathbb{E}\left[h\left(\varepsilon_{i}, \varepsilon_{j}\right)\right]\right\}\right| .
$$

For any $\gamma \geq 0$, Markov's inequality implies

$$
\mathbb{P}\left(\mathcal{Z}_{M}>t\right)=\mathbb{P}\left(e^{\gamma \mathcal{Z}_{M}}>e^{\gamma t}\right) \leq e^{-\gamma t} \mathbb{E} e^{\gamma \mathcal{Z}_{M}}
$$

Let $M_{n}=\lfloor n / 2\rfloor$, the largest integer that is no larger than $n / 2$. Let $\sigma_{i}, i=1, \ldots, n$ denote a Rademacher sequence independent of $\left\{\mathbf{X}_{i}, \varepsilon_{i}\right\}_{i=1}^{n}$. In the following, we denote $\Gamma=\{\mathbf{A} \in$ $\left.\mathbb{R}^{m_{1} \times m_{2}} \mid\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1} \leq M\right\}$. By Lemma A. 1 of Clémençon et al. (2008) and the convexity of $\exp (x)$, we have

$$
\mathbb{E} e^{\gamma \mathcal{Z}_{M}} \leq \frac{1}{n!} \sum_{\pi} \mathbb{E} \exp \left(\gamma \sup _{\mathbf{A} \in \Gamma}\left|M_{n}^{-1} \sum_{i=1}^{M_{n}}\left\{h\left(\varepsilon_{\pi(i)}, \varepsilon_{\pi\left(M_{n}+i\right)}\right)-\mathbb{E} h\left(\varepsilon_{\pi(i)}, \varepsilon_{\pi\left(M_{n}+i\right)}\right)\right\}\right|\right)
$$

where the first summation taking over all permutations $\pi$ of $\{1, \ldots, n\}$. Applying the symmetrization theorem and the contraction theorem, we obtain

$$
\mathbb{E} e^{\gamma \mathcal{Z}_{M}} \leq \frac{1}{n!} \sum_{\pi} \mathbb{E} \exp \left(4 \gamma \sup _{\mathbf{A} \in \Gamma}\left|M_{n}^{-1} \sum_{i=1}^{M_{n}} \sigma_{i}\left\langle\mathbf{X}_{\pi(i)}-\mathbf{X}_{\pi\left(M_{n}+i\right)}, \mathbf{A}-\mathbf{A}_{0}\right\rangle\right|\right) .
$$

Then by the dual norm inequality we have

$$
\begin{aligned}
\mathbb{E} e^{\gamma \mathcal{Z}_{M}} & \leq \frac{1}{n!} \sum_{\pi} \mathbb{E} \exp \left(\begin{array}{r}
\left.4 \gamma \sup _{\mathbf{A} \in \Gamma}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1}\left\|M_{n}^{-1} \sum_{i=1}^{M_{n}} \sigma_{i}\left(\mathbf{X}_{\pi(i)}-\mathbf{X}_{\pi\left(M_{n}+i\right)}\right)\right\|_{\mathrm{op}}\right) \\
\\
\end{array} \leq \frac{1}{n!} \sum_{\pi} \mathbb{E} \exp \left(4 \gamma M\left\|M_{n}^{-1} \sum_{i=1}^{M_{n}} \sigma_{i}\left(\mathbf{X}_{\pi(i)}-\mathbf{X}_{\pi\left(M_{n}+i\right)}\right)\right\|_{\mathrm{op}}\right)\right.
\end{aligned}
$$

Next we adopt the $\varepsilon$-net argument in the Chapter 4 of Vershynin (2018) to bound the last term. Choose $\varepsilon=1 / 4$, by Corollary 4.2.13 of Vershynin (2018), we can find an $\varepsilon$-net $\mathcal{N}$ of the sphere $\mathcal{D}^{m_{1}-1}$ and $\varepsilon$-net $\mathcal{M}$ of the sphere $\mathcal{D}^{m_{2}-1}$ with cardinalities $|\mathcal{N}| \leq 9^{m_{1}}$ and $|\mathcal{M}| \leq 9^{m_{2}}$ such that for any $\mathbf{A} \in \mathbb{R}^{m_{1} \times m_{2}}$, the operator norm of $\mathbf{A}$ can be bounded using these nets as follows: $\|\mathbf{A}\|_{\text {op }} \leq 2 \max _{u \in \mathcal{N}, v \in \mathcal{M}} u^{\top} \mathbf{A} v$. Then we have

$$
\begin{aligned}
& \mathbb{E} \exp \left(4 \gamma M\left\|M_{n}^{-1} \sum_{i=1}^{M_{n}} \sigma_{i}\left(\mathbf{X}_{\pi(i)}-\mathbf{X}_{\pi\left(M_{n}+i\right)}\right)\right\|_{\mathrm{op}}\right) \\
\leq & \mathbb{E} \exp \left(8 \gamma M \max _{u \in \mathcal{N}, v \in \mathcal{M}} u^{\top}\left\{M_{n}^{-1} \sum_{i=1}^{M_{n}} \sigma_{i}\left(\mathbf{X}_{\pi(i)}-\mathbf{X}_{\pi\left(M_{n}+i\right)}\right)\right\} v\right) \\
\leq & 2 \cdot 9^{m_{1}+m_{2}} \max _{u \in \mathcal{N}, v \in \mathcal{M}} \mathbb{E} \exp \left(8 \gamma M u^{\top}\left\{M_{n}^{-1} \sum_{i=1}^{M_{n}} \sigma_{i}\left(\mathbf{X}_{\pi(i)}-\mathbf{X}_{\pi\left(M_{n}+i\right)}\right)\right\} v\right) \\
= & 2 \cdot 9^{m_{1}+m_{2}} \max _{u \in \mathcal{N}, v \in \mathcal{M}} \mathbb{E} \exp \left(8 \gamma M M_{n}^{-1} \sum_{i=1}^{M_{n}} u^{\top}\left(\mathbf{X}_{\pi(i)}-\mathbf{X}_{\pi\left(M_{n}+i\right)}\right) v\right) \\
\leq & 2 \cdot 9^{m_{1}+m_{2}} \exp \left(\frac{C_{1}^{2} \kappa_{0}^{2} \gamma^{2} M^{2}}{n}\right),
\end{aligned}
$$

where the second inequality follows from the simple bound

$$
\mathbb{E}\left[\max _{1 \leq j \leq p} e^{\left|z_{j}\right|}\right] \leq p \max _{1 \leq j \leq p} \mathbb{E}\left[e^{\left|z_{j}\right|}\right] \leq p \max _{1 \leq j \leq p} \mathbb{E}\left[e^{z_{j}}+e^{-z_{j}}\right] \leq 2 p \max _{1 \leq j \leq p} \mathbb{E}\left[e^{z_{j}}\right]
$$

holding for symmetric random variables $z_{j}$, and the last inequality follows the sub-Gaussian assumption on the covariate. Above discussion implies

$$
\begin{aligned}
\mathbb{P}\left(\mathcal{Z}_{M}>t\right) & \leq \inf _{\gamma \geq 0} e^{-\gamma t} \mathbb{E} e^{\gamma \mathcal{Z}_{M}} \\
& \leq \inf _{\gamma \geq 0} 2 \cdot 9^{m_{1}+m_{2}} \exp \left(C_{1}^{2} \kappa_{0}^{2} \gamma^{2} M^{2} / n-\gamma t\right) \\
& =2 \cdot 9^{m_{1}+m_{2}} \exp \left(-\frac{n t^{2}}{4 C_{1}^{2} \kappa_{0}^{2} M^{2}}\right) .
\end{aligned}
$$

Let $t=8 C_{1} \kappa_{0} \sqrt{\frac{m_{1}+m_{2}}{n}} M$, then we conclude that $\mathcal{Z}_{M} \leq 8 C_{1} \kappa_{0} \sqrt{\frac{m_{1}+m_{2}}{n}} M$ with probability at least $1-2 \exp \left(-4\left(m_{1}+m_{2}\right)\right)$.

For the second conclusion, to obtain a uniform bound for all $\mathbf{A} \in \mathcal{S}$, we apply the first conclusion and peeling device given in Van de Geer (2000) and Elsener and van de Geer (2018). Without loss of generality we assume $m_{1} \geq m_{2}$, then we can subdivide the set $\mathcal{S}$ as follows

$$
\mathcal{S}=\left\{\mathbf{A} \in \mathcal{S}:\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1} \leq \frac{m_{1}+m_{2}}{n}\right\} \bigcup\left\{\mathbf{A} \in \mathcal{S}: \frac{m_{1}+m_{2}}{n}<\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1} \leq 2 \sqrt{m_{1}^{2} m_{2}} \eta\right\}
$$

On the first set $\mathcal{S}_{0}=\left\{\mathbf{A} \in \mathcal{S}:\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1} \leq \frac{m_{1}+m_{2}}{n}\right\}$

$$
\begin{aligned}
& \mathbb{P}\left(\exists \mathbf{A} \in \mathcal{S}_{0}:\left|\left\{Q_{n}(\mathbf{A})-Q(\mathbf{A})\right\}-\left\{Q_{n}\left(\mathbf{A}_{0}\right)-Q\left(\mathbf{A}_{0}\right)\right\}\right|>\lambda_{\varepsilon}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1}+\widetilde{\lambda}\right) \\
& \leq \mathbb{P}\left(\exists \mathbf{A} \in \mathcal{S}_{0}:\left|\left\{Q_{n}(\mathbf{A})-Q(\mathbf{A})\right\}-\left\{Q_{n}\left(\mathbf{A}_{0}\right)-Q\left(\mathbf{A}_{0}\right)\right\}\right|>\frac{m_{1}+m_{2}}{n} \lambda_{\varepsilon}^{\prime}\right) \\
& \leq 2 \exp \left(-4\left(m_{1}+m_{2}\right)\right) .
\end{aligned}
$$

We further subdivide the second set $\left\{\mathbf{A} \in \mathcal{S}: \frac{m_{1}+m_{2}}{n}<\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1} \leq 2 \sqrt{m_{1}^{2} m_{2}} \eta\right\}$ into

$$
\begin{aligned}
& \left\{\mathbf{A} \in \mathcal{S}: \frac{m_{1}+m_{2}}{n}<\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1} \leq 2 \sqrt{m_{1}^{2} m_{2}} \eta\right\} \\
\subset & \bigcup_{k=1}^{k_{0}}\left\{\mathbf{A} \in \mathcal{S}: \frac{m_{1}+m_{2}}{n} 2^{k}<\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1} \leq \frac{m_{1}+m_{2}}{n} 2^{k+1}\right\}:=\bigcup_{k=1}^{k_{0}} \mathcal{S}_{k},
\end{aligned}
$$

where $k_{0}$ is the smallest integer such that $\frac{m_{1}+m_{2}}{n} 2^{k_{0}+1} \geq 2 \sqrt{m_{1}^{2} m_{2}} \eta$. For each $\mathcal{S}_{k}$

$$
\begin{aligned}
& \mathbb{P}\left(\exists \mathbf{A} \in \mathcal{S}_{k}:\left|\left\{Q_{n}(\mathbf{A})-Q(\mathbf{A})\right\}-\left\{Q_{n}\left(\mathbf{A}_{0}\right)-Q\left(\mathbf{A}_{0}\right)\right\}\right|>\lambda_{\varepsilon}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1}+\widetilde{\lambda}\right) \\
& \leq \mathbb{P}\left(\exists \mathbf{A} \in \mathcal{S}_{k}:\left|\left\{Q_{n}(\mathbf{A})-Q(\mathbf{A})\right\}-\left\{Q_{n}\left(\mathbf{A}_{0}\right)-Q\left(\mathbf{A}_{0}\right)\right\}\right|>\frac{m_{1}+m_{2}}{n} 2^{k+1} \lambda_{\varepsilon}^{\prime}\right) \\
& \leq 2 \exp \left(-4\left(m_{1}+m_{2}\right)\right) .
\end{aligned}
$$

Note that $k_{0}=\log _{2}\left(\frac{n \sqrt{m_{1}^{2} m_{2}} \eta}{m_{1}+m_{2}}\right)$, by the high-dimensional setting $n<m_{1} m_{2}$, there exists a constant $c$ such that $\left(k_{0}+1\right) \exp \left(-2\left(m_{1}+m_{2}\right)\right) \leq 1$ when $m_{1}+m_{2}>c$. Then by the
union bound

$$
\begin{aligned}
& \mathbb{P}\left(\exists \mathbf{A} \in \mathcal{S}:\left|\left\{Q_{n}(\mathbf{A})-Q(\mathbf{A})\right\}-\left\{Q_{n}\left(\mathbf{A}_{0}\right)-Q\left(\mathbf{A}_{0}\right)\right\}\right|>\lambda_{\varepsilon}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1}+\widetilde{\lambda}\right) \\
& \leq \sum_{k=0}^{k_{0}} \mathbb{P}\left(\exists \mathbf{A} \in \mathcal{S}_{k}:\left|\left\{Q_{n}(\mathbf{A})-Q(\mathbf{A})\right\}-\left\{Q_{n}\left(\mathbf{A}_{0}\right)-Q\left(\mathbf{A}_{0}\right)\right\}\right|>\lambda_{\varepsilon}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1}+\widetilde{\lambda}\right) \\
& \leq 2\left(k_{0}+1\right) \exp \left(-4\left(m_{1}+m_{2}\right)\right) \leq 2 \exp \left(-2\left(m_{1}+m_{2}\right)\right)
\end{aligned}
$$

which completes the proof of the second conclusion.
(b) The first conclusion follows from the definition of $\lambda^{*}$. Next, we give the proof of the second conclusion. Conditional on $\left\{\mathbf{X}_{i}\right\}_{i=1}^{n}$. For any $\gamma \geq 0$, Markov's inequality implies

$$
\mathbb{P}\left(S_{n}>t \mid \mathbf{X}\right)=\mathbb{P}\left(e^{\gamma S_{n}}>e^{\gamma t} \mid \mathbf{X}\right) \leq e^{-\gamma t} \mathbb{E}_{\mathbb{P} \mid \mathbf{X}} e^{\gamma S_{n}}
$$

Recall that $S_{n}=\{n(n-1)\}^{-1}\left\|\sum_{i \neq j}\left(\mathbf{X}_{j}-\mathbf{X}_{i}\right) \operatorname{sign}\left(\varepsilon_{i}-\varepsilon_{j}\right)\right\|_{\mathrm{op}}$, then using the $\varepsilon$-net, we have

$$
\begin{aligned}
& \mathbb{E}_{\mathbb{P} \mid \mathbf{X}} e^{\gamma S_{n}} \\
\leq & 2 \cdot 9^{m_{1}+m_{2}} \max _{u \in \mathcal{N}, v \in \mathcal{M}} \mathbb{E}_{\mathbb{P} \mid \mathbf{X}} \exp \left(2 \gamma\{n(n-1)\}^{-1} \sum_{i \neq j} u^{\top}\left(\mathbf{X}_{j}-\mathbf{X}_{i}\right) v \operatorname{sign}\left(\varepsilon_{i}-\varepsilon_{j}\right)\right) \\
\leq & 2 \cdot 9^{m_{1}+m_{2}} \max _{u \in \mathcal{N}, v \in \mathcal{M}} \mathbb{E}_{\mathbb{P} \mid \mathbf{X}} \exp \left(2 \gamma \frac{1}{n!} \sum_{\pi} M_{n}^{-1} \sum_{i=1}^{M_{n}} u^{\top}\left(\mathbf{X}_{\pi(i)}-\mathbf{X}_{\pi\left(M_{n}+i\right)}\right) v \operatorname{sign}\left(\varepsilon_{\pi(i)}-\varepsilon_{\pi\left(M_{n}+i\right)}\right)\right) \\
\leq & \frac{1}{n!} \sum_{\pi} 2 \cdot 9^{m_{1}+m_{2}} \max _{u \in \mathcal{N}, v \in \mathcal{M}} \mathbb{E}_{\mathbb{P} \mid \mathbf{X}} \exp \left(2 \gamma M_{n}^{-1} \sum_{i=1}^{M_{n}} u^{\top}\left(\mathbf{X}_{\pi(i)}-\mathbf{X}_{\pi\left(M_{n}+i\right)}\right) v \operatorname{sign}\left(\varepsilon_{\pi(i)}-\varepsilon_{\pi\left(M_{n}+i\right)}\right)\right) \\
\leq & \frac{1}{n!} \sum_{\pi} 2 \cdot 9^{m_{1}+m_{2}} \max _{u \in \mathcal{N}, v \in \mathcal{M}} \exp \left(4 \frac{\gamma^{2}}{M_{n}^{2}} \sum_{i=1}^{M_{n}} \frac{\left\{u^{\top}\left(\mathbf{X}_{\pi(i)}-\mathbf{X}_{\pi\left(M_{n}+i\right)}\right) v\right\}^{2}}{2}\right) \\
\leq & \frac{1}{n!} \sum_{\pi} 2 \cdot 9^{m_{1}+m_{2}} \exp \left(16 \frac{\gamma^{2}}{n} \max _{u \in \mathcal{N}, v \in \mathcal{M}} \frac{1}{n} \sum_{i=1}^{n}\left(u^{\top} \mathbf{X}_{i} v\right)^{2}\right) .
\end{aligned}
$$

Assume $16 \max _{u \in \mathcal{N}, v \in \mathcal{M}} \frac{1}{n} \sum_{i=1}^{n}\left(u^{\top} \mathbf{X}_{i} v\right)^{2} \leq \kappa_{0}^{2} C_{2}^{2}$ for some constant $C_{2}$, we obtain

$$
\begin{aligned}
\mathbb{P}\left(S_{n}>t \mid \mathbf{X}\right) & \leq \inf _{\gamma \geq 0} e^{-\gamma t} \mathbb{E}_{\mathbb{P} \mid \mathbf{X}} e^{\gamma S_{n}} \\
& \leq \inf _{\gamma \geq 0} 2 \cdot 9^{m_{1}+m_{2}} \exp \left(C_{2}^{2} \kappa_{0}^{2} \gamma^{2} / n-\gamma t\right) \\
& =2 \cdot 9^{m_{1}+m_{2}} \exp \left(-\frac{n t^{2}}{4 C_{2}^{2} \kappa_{0}^{2}}\right)
\end{aligned}
$$

Let $t=4 C_{2} \kappa_{0} \sqrt{\frac{m_{1}+m_{2}}{n}}$, it follows that $S_{n} \leq 4 C_{2} \kappa_{0} \sqrt{\frac{m_{1}+m_{2}}{n}}$ with probability at least $1-$ $2 \exp \left(-\left(m_{1}+m_{2}\right)\right)$. Then if $m_{1}+m_{2}>\ln \left(2 / \alpha_{0}\right)$, we conclude that $\lambda^{*}<4 C_{2} \kappa_{0} \sqrt{\frac{m_{1}+m_{2}}{n}}$.

Finally we only need to show that there exists some constant $C_{2}$ such that

$$
16 \max _{u \in \mathcal{N}, v \in \mathcal{M}} \frac{1}{n} \sum_{i=1}^{n}\left(u^{\top} \mathbf{X}_{i} v\right)^{2} \leq \kappa_{0}^{2} C_{2}^{2}
$$

holds with high probability. Under the Assumption 12, it's easy to show that $u^{\top} \mathbf{X}_{i} v / \kappa_{0}$ is sub-Gaussian with $\left\|u^{\top} \mathbf{X}_{i} v / \kappa_{0}\right\|_{\psi_{2}} \leq 1$ and $\mathbb{E}\left(u^{\top} \mathbf{X}_{i} v / \kappa_{0}\right)^{2} \leq \mu$ for all $u \in \mathcal{N}, v \in \mathcal{M}$, where $\mu$ is a constant. By $\left\|\left(u^{\top} \mathbf{X}_{i} v / \kappa_{0}\right)^{2}\right\|_{\psi_{1}}=\left(\left\|u^{\top} \mathbf{X}_{i} v / \kappa_{0}\right\|_{\psi_{2}}\right)^{2}$ we have that $\left(u^{\top} \mathbf{X}_{i} v / \kappa_{0}\right)^{2}$ is sub-exponential with constant parameter ( $\nu, \alpha$ ) and

$$
\mathbb{P}\left(\frac{1}{n} \sum_{i=1}^{n}\left(\left(u^{\top} \mathbf{X}_{i} v / \kappa_{0}\right)^{2}-\mathbb{E}\left(u^{\top} \mathbf{X}_{i} v / \kappa_{0}\right)^{2}\right)>t\right) \leq\left\{\begin{array}{ll}
e^{-\frac{n t^{2}}{2 \nu^{2}}} & \text { for } 0 \leq t \leq \frac{\nu^{2}}{\alpha} \\
e^{-\frac{n t}{2 \alpha}} & \text { for } t>\frac{\nu^{2}}{\alpha}
\end{array} .\right.
$$

We take $t=\frac{\nu^{2}}{\alpha}$, then

$$
\begin{aligned}
\mathbb{P}\left(\max _{u \in \mathcal{N}, v \in \mathcal{M}} \frac{1}{n} \sum_{i=1}^{n}\left(u^{\top} \mathbf{X}_{i} v / \kappa_{0}\right)^{2}>\mu+\frac{\nu^{2}}{\alpha}\right) & \leq \sum_{k=1}^{9^{m_{1}+m_{2}}} \mathbb{P}\left(\frac{1}{n} \sum_{i=1}^{n}\left(u^{\top} \mathbf{X}_{i} v / \kappa_{0}\right)^{2}>\mu+\frac{\nu^{2}}{\alpha}\right) \\
& \leq 9^{m_{1}+m_{2}} \exp \left(-\frac{\nu^{2}}{2 \alpha^{2}} n\right) .
\end{aligned}
$$

Hence there exists constant $C_{2}$ and $c_{2}$ such that

$$
\mathbb{P}\left(16 \max _{u \in \mathcal{N}, v \in \mathcal{M}} \frac{1}{n} \sum_{i=1}^{n}\left(u^{\top} \mathbf{X}_{i} v\right)^{2} \leq \kappa_{0}^{2} C_{2}^{2}\right) \geq 1-9^{m_{1}+m_{2}} \exp \left(-c_{2} n\right) .
$$

Note that $16 \max _{u \in \mathcal{N}, v \in \mathcal{M}} \frac{1}{n} \sum_{i=1}^{n}\left(u^{\top} \mathbf{X}_{i} v\right)^{2} \leq \kappa_{0}^{2} C_{2}^{2}$ with probability at least $1-9^{m_{1}+m_{2}} \exp \left(-c_{2} n\right)$. If $n>4\left(m_{1}+m_{2}\right) / c_{2}$, we obtain $\lambda^{*}<4 C_{2} \kappa_{0} \sqrt{\frac{m_{1}+m_{2}}{n}}$ with probability at least 1 $\exp \left(-\left(m_{1}+m_{2}\right)\right)$.
(c) We denote $\operatorname{Cov}\left(\operatorname{vec}\left(\mathbf{X}_{1}\right)\right)$ by $\mathbf{J}$. First we show that, for any $\mathbf{A}$

$$
Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right) \geq \frac{b_{1}}{2}\left\|\mathbf{J}^{\frac{1}{2}} \operatorname{vec}\left(\mathbf{A}-\mathbf{A}_{0}\right)\right\|_{2}^{2} \wedge \frac{b_{1} b_{4}}{2}\left\|\mathbf{J}^{\frac{1}{2}} \operatorname{vec}\left(\mathbf{A}-\mathbf{A}_{0}\right)\right\|_{2}
$$

which implies that $\mathbf{A}_{0}$ is the minimizer of $Q(\mathbf{A})$. We define the maximal radius over which the criterion function can be minorated by a quadratic function

$$
r_{\mathbf{A}_{0}}=\sup _{r}\left\{r: Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right) \geq \frac{b_{1}}{2}\left\|\mathbf{J}^{1 / 2} \operatorname{vec}\left(\mathbf{A}-\mathbf{A}_{0}\right)\right\|_{2}^{2} \text { for all }\left\|\mathbf{J}^{1 / 2} \operatorname{vec}\left(\mathbf{A}-\mathbf{A}_{0}\right)\right\|_{2} \leq r\right\}
$$

We claim that $r_{\mathbf{A}_{0}} \geq b_{4}$. By Knight's identity (Koenker, 2005),

$$
\begin{aligned}
Q_{n}(\mathbf{A})-Q_{n}\left(\mathbf{A}_{0}\right)= & \{n(n-1)\}^{-1} \sum_{i \neq j}\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle\left\{\mathrm{I}\left(\zeta_{i j}<0\right)-1 / 2\right\} \\
& +\{n(n-1)\}^{-1} \sum_{i \neq j} \int_{0}^{\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle}\left\{\mathrm{I}\left(\zeta_{i j} \leq s\right)-\mathrm{I}\left(\zeta_{i j} \leq 0\right)\right\} d s .
\end{aligned}
$$

Denote the distribution function of $\zeta_{i j}$ by $F^{*}(\cdot)$. By the independence of $\mathbf{X}$ and $\varepsilon$, we have

$$
Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right)=\mathbb{E}_{\mathbf{X}} \int_{0}^{\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle}\left\{F^{*}(s)-F^{*}(0)\right\} d s
$$

By the Taylor expansion and Assumption 2, for some $\xi_{i j}$ between 0 and $\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle$, we have

$$
\begin{aligned}
Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right) & =\mathbb{E}_{\mathbf{X}} \int_{0}^{\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle}\left\{f^{*}(0) s+\frac{\partial f^{*}\left(\xi_{i j}\right)}{\partial t} \frac{s^{2}}{2}\right\} d s \\
& \geq \frac{b_{1}}{2} \mathbb{E}_{\mathbf{X}}\left[\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle^{2}\right]-\frac{b_{2}}{6} \mathbb{E}_{\mathbf{X}}\left[\left|\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle\right|^{3}\right]
\end{aligned}
$$

Note that for all $\mathbf{A}-\mathbf{A}_{0}$, if

$$
\left\|\mathbf{J}^{1 / 2} \operatorname{vec}\left(\mathbf{A}-\mathbf{A}_{0}\right)\right\|_{2} \leq b_{4} \leq \frac{3 b_{1}}{2 \sqrt{2} b_{2}} \inf _{\mathbf{A}-\mathbf{A}_{0} \neq 0} \frac{\left(\mathbb{E}\left\langle\mathbf{X}_{1}-\mathbf{X}_{2}, \mathbf{A}-\mathbf{A}_{0}\right\rangle^{2}\right)^{3 / 2}}{\mathbb{E}\left|\left\langle\mathbf{X}_{1}-\mathbf{X}_{2}, \mathbf{A}-\mathbf{A}_{0}\right\rangle\right|^{3}}
$$

it follows that

$$
\frac{b_{1}}{4} \mathbb{E}_{\mathbf{X}}\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle^{2} \geq \frac{b_{2}}{6} \mathbb{E}_{\mathbf{X}}\left|\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle\right|^{3}
$$

Hence for any $\left\|\mathbf{J}^{1 / 2} \operatorname{vec}\left(\mathbf{A}-\mathbf{A}_{0}\right)\right\|_{2} \leq b_{4}$, we have

$$
\begin{aligned}
Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right) & \geq \frac{b_{1}}{4} \mathbb{E}_{\mathbf{X}}\left(\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle^{2}\right) \\
& =\frac{b_{1}}{2}\left\|\mathbf{J}^{\frac{1}{2}} \operatorname{vec}\left(\mathbf{A}-\mathbf{A}_{0}\right)\right\|_{2}^{2}
\end{aligned}
$$

This implies $r_{\mathbf{A}_{0}} \geq b_{4}$.
By construction of $r_{\mathbf{A}_{0}}$ and the convexity of $Q$, for any $\mathbf{A}$ we have

$$
\begin{aligned}
& Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right) \\
\geq & \frac{b_{1}}{2}\left\|\mathbf{J}^{1 / 2} \operatorname{vec}\left(\mathbf{A}-\mathbf{A}_{0}\right)\right\|_{2}^{2} \wedge\left\{\frac{\left\|\mathbf{J}^{1 / 2} \operatorname{vec}\left(\mathbf{A}-\mathbf{A}_{0}\right)\right\|_{2}}{r_{\mathbf{A}_{0}}} \cdot \underset{\left\|\mathbf{J}^{1 / 2} \operatorname{vec}\left(\tilde{\mathbf{A}}-\mathbf{A}_{0}\right)\right\|_{2}=r_{\mathbf{A}_{0}}}{ } Q(\widetilde{\mathbf{A}})-Q\left(\mathbf{A}_{0}\right)\right\} \\
\geq & \frac{b_{1}}{2}\left\|\mathbf{J}^{1 / 2} \operatorname{vec}\left(\mathbf{A}-\mathbf{A}_{0}\right)\right\|_{2}^{2} \wedge\left\{\frac{\left\|\mathbf{J}^{1 / 2} \operatorname{vec}\left(\mathbf{A}-\mathbf{A}_{0}\right)\right\|_{2}}{r_{\mathbf{A}_{0}}} \frac{b_{1}}{2} r_{\mathbf{A}_{0}}^{2}\right\} \\
\geq & \frac{b_{1}}{2}\left\|\mathbf{J}^{1 / 2} \operatorname{vec}\left(\mathbf{A}-\mathbf{A}_{0}\right)\right\|_{2}^{2} \wedge\left\{\frac{b_{1} b_{4}}{2}\left\|\mathbf{J}^{1 / 2} \operatorname{vec}\left(\mathbf{A}-\mathbf{A}_{0}\right)\right\|_{2}\right\}
\end{aligned}
$$

which implies $\mathbf{A}_{0}$ is the minimizer of $Q(\mathbf{A})$. By Assumption 3, we further have

$$
Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right) \geq \frac{b_{1} b_{3}}{2}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{F}^{2} \wedge \frac{b_{1} b_{3} b_{4}}{2}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{F}
$$

for all $\mathbf{A}-\mathbf{A}_{0} \in \mathcal{C}$. Using the fact that $\widehat{\mathbf{A}}$ is the minimizer of the objective function:

$$
Q_{n}(\widehat{\mathbf{A}})+\lambda\|\widehat{\mathbf{A}}\|_{1} \leq Q_{n}\left(\mathbf{A}_{0}\right)+\lambda\left\|\mathbf{A}_{0}\right\|_{1}
$$

Similar to the proof of Theorem 1, we obtain

$$
Q(\widehat{\mathbf{A}})-Q\left(\mathbf{A}_{0}\right) \leq\left(\lambda_{\varepsilon}+\lambda\right)\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{1}+\widetilde{\lambda}
$$

By Lemma 21, we have $\widehat{\mathbf{A}}-\mathbf{A}_{0} \in \mathcal{C}$, then

$$
\begin{aligned}
\frac{b_{1} b_{3}}{2}\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2} \wedge \frac{b_{1} b_{3} b_{4}}{2}\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F} & \leq\left(\lambda_{\varepsilon}+\lambda\right)\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{1}+\widetilde{\lambda} \\
& \leq\left(\lambda_{\varepsilon}+\lambda\right)\left\{4 \sqrt{2 r}\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}+4 \sum_{j=r+1}^{m} \sigma_{j}\left(\mathbf{A}_{0}\right)\right\}+\widetilde{\lambda} \\
& \leq\left(\lambda_{\varepsilon}+\lambda\right) \max \left\{8 \sqrt{2 r}\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}, 8 \sum_{j=r+1}^{m} \sigma_{j}\left(\mathbf{A}_{0}\right)\right\}+\widetilde{\lambda} .
\end{aligned}
$$

If $\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F} \leq b_{4}$, combining the fact that $\sqrt{\tilde{\lambda}}$ is dominated by $\lambda_{\varepsilon}$ when $n \gtrsim\left(m_{1}+m_{2}\right)$, it follows that $Q(\widehat{\mathbf{A}})-Q\left(\mathbf{A}_{0}\right) \geq \frac{b_{1} b_{3}}{2}\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2}$ and

$$
\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F} \leq \max \left\{\frac{32\left(\lambda+\lambda_{\varepsilon}\right) \sqrt{r}}{b_{1} b_{3}},\left\{\frac{16\left(\lambda+\lambda_{\varepsilon}\right) \sum_{j=r+1}^{m} \sigma_{j}\left(\mathbf{A}_{0}\right)}{b_{1} b_{3}}\right\}^{1 / 2}\right\}
$$

If $\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F} \geq b_{4}$, then

$$
\frac{b_{1} b_{3} b_{4}}{2}\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F} \leq 2\left(\lambda_{\varepsilon}+\lambda\right)\left\{4 \sqrt{2 r}\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}+4 \sum_{j=r+1}^{m} \sigma_{j}\left(\mathbf{A}_{0}\right)\right\}
$$

which implies

$$
1 \leq \frac{4\left(\lambda_{\varepsilon}+\lambda\right)}{b_{1} b_{3} b_{4}}\left\{4 \sqrt{2 r}+4 \frac{\sum_{j=r+1}^{m} \sigma_{j}\left(\mathbf{A}_{0}\right)}{\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}}\right\} \leq \frac{4\left(\lambda_{\varepsilon}+\lambda\right)}{b_{1} b_{3} b_{4}}\left\{4 \sqrt{2 r}+4 \frac{\sum_{j=r+1}^{m} \sigma_{j}\left(\mathbf{A}_{0}\right)}{b_{4}}\right\}
$$

Similar to the proof of Theorem 1, we choose $r=\#\left\{j \in\{1,2, \ldots, p\} \mid \sigma_{j}\left(\mathbf{A}_{0}\right) \geq \tau\right\}$, then we have $\sum_{j=r+1}^{p} \sigma_{j}\left(\mathbf{A}_{0}\right) \leq \tau^{1-q} R_{q}$ and $r \leq R_{q} \tau^{-q}$. Taking $\kappa=b_{1} b_{3} / 2$ and $\tau=\left(\lambda+\lambda_{\varepsilon}\right) / \kappa$, this gives

$$
1 \leq\left\{\sqrt{2}+\frac{1}{b_{4}}\right\} \frac{16}{b_{1} b_{3} b_{4}}\left(\frac{\lambda+\lambda_{\varepsilon}}{\kappa}\right)^{1-q / 2} R_{q}^{1 / 2}<1
$$

which is a contradiction. Hence if $\frac{16 \sqrt{2} b_{4}+16}{b_{1} b_{3} b_{4}^{2}}\left(\frac{\lambda+\lambda_{\varepsilon}}{\kappa}\right)^{1-q / 2} R_{q}^{1 / 2}<1$, we have

$$
Q(\widehat{\mathbf{A}})-Q\left(\mathbf{A}_{0}\right) \geq \frac{b_{1} b_{3}}{2}\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2}
$$

## Appendix C: Proof of Corollary 10

In the proof of Theorem 6, we use $\varepsilon$-net argument to bound $\left\|M_{n}^{-1} \sum_{i=1}^{M_{n}} \sigma_{i}\left(\mathbf{X}_{\pi(i)}-\mathbf{X}_{\pi\left(M_{n}+i\right)}\right)\right\|_{\text {op }}$ and $\left\|\sum_{i \neq j}\left(\mathbf{X}_{i}-\mathbf{X}_{j}\right) \operatorname{sign}\left(\varepsilon_{i}-\varepsilon_{j}\right)\right\|_{\mathrm{op}}$. But in the high dimensional linear regression model, due to the diagonal structure of $\mathbf{X}$, we bound them in the following way to get a faster rate:

$$
\begin{gathered}
\left\|M_{n}^{-1} \sum_{i=1}^{M_{n}} \sigma_{i}\left(\mathbf{X}_{\pi(i)}-\mathbf{X}_{\pi\left(M_{n}+i\right)}\right)\right\|_{\mathrm{op}} \leq \max _{1 \leq j \leq d}\left|M_{n}^{-1} \sum_{i=1}^{M_{n}} \sigma_{i}\left(\mathbf{X}_{\pi(i) j j}-\mathbf{X}_{\pi\left(M_{n}+i\right) j j}\right)\right|, \\
\left\|\sum_{i \neq j}\left(\mathbf{X}_{i}-\mathbf{X}_{j}\right) \operatorname{sign}\left(\varepsilon_{i}-\varepsilon_{j}\right)\right\|_{\mathrm{op}} \leq \max _{1 \leq k \leq d}\left|\sum_{i \neq j}\left(\mathbf{X}_{i k k}-\mathbf{X}_{j k k}\right) \operatorname{sign}\left(\varepsilon_{i}-\varepsilon_{j}\right)\right|,
\end{gathered}
$$

where $\mathbf{X}_{\pi(i) j j}$ and $\mathbf{X}_{i k k}$ denote the $j$-th and $k$-th diagonal element of $\mathbf{X}_{\pi(i)}$ and $\mathbf{X}_{i}$ respectively. Then with the two inequalities, similar to the proof of Theorem 6 , we can get the conclusion in Corollary 10.

## Appendix D: Proof of Theorem 14

Combining the following lemma with Theorem 1, we can complete the proof of Theorem 14.

Lemma 24 Suppose that Assumptions 11-13 hold. Suppose $\mathbf{A}_{0} \in \mathcal{B}_{q}\left(R_{q}\right) \cap \mathcal{S}$. Then we have the following conclusions.
(a) Define for all $M>0$

$$
\mathcal{Z}_{M}:=\sup _{\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1} \leq M}\left|\left\{Q_{n}(\mathbf{A})-Q(\mathbf{A})\right\}-\left\{Q_{n}\left(\mathbf{A}_{0}\right)-Q\left(\mathbf{A}_{0}\right)\right\}\right| .
$$

Then there exists a universal constant $C_{1}>0$, with $\lambda_{\varepsilon}^{\prime}=8 C_{1} \kappa_{0} \sqrt{\frac{m_{1}+m_{2}}{n}}$, such that

$$
\mathcal{Z}_{M} \leq \lambda_{\varepsilon}^{\prime} M
$$

with probability at least $1-2 \exp \left(-4\left(m_{1}+m_{2}\right)\right)$. Moreover, let $\lambda_{\varepsilon}=2 \lambda_{\varepsilon}^{\prime}$ and $\widetilde{\lambda}=$ $2 \frac{m_{1}+m_{2}}{n} \lambda_{\varepsilon}^{\prime}$, we have, for any $\mathbf{A} \in \mathcal{S}$

$$
\left|\left\{Q_{n}(\mathbf{A})-Q(\mathbf{A})\right\}-\left\{Q_{n}\left(\mathbf{A}_{0}\right)-Q\left(\mathbf{A}_{0}\right)\right\}\right| \leq \lambda_{\varepsilon}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1}+\widetilde{\lambda}
$$

with probability at least $1-2 \exp \left(-2\left(m_{1}+m_{2}\right)\right)$.
(b) $\lambda^{*}>2\left\|\nabla Q_{n}\left(\mathbf{A}_{0}\right)\right\|_{\text {op }}$ with probability $1-\alpha_{0}$. If $m_{1}+m_{2}>\ln \left(2 / \alpha_{0}\right)$, then there exists constant $C_{2}>0$ such that $\lambda^{*} \leq 4 C_{2} \kappa_{0} \sqrt{\frac{m_{1}+m_{2}}{n}}$ with probability at least $1-$ $\exp \left(-\left(m_{1}+m_{2}\right)\right)$.
(c) $\mathbf{A}_{0}$ is the minimizer of $Q(\mathbf{A})$. Take $\kappa=b_{1} b_{3} / 2$, if $\frac{16 \sqrt{2} b_{4}+16}{b_{1} b_{3} b_{4}^{2}}\left(\frac{\lambda+\lambda_{\varepsilon}}{\kappa}\right)^{1-q / 2} R_{q}^{1 / 2}<1$, then we have

$$
Q(\widehat{\mathbf{A}})-Q\left(\mathbf{A}_{0}\right) \geq \frac{b_{1} b_{3}}{2}\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2}
$$

Proof [Proof of Lemma 24] (a) We only prove the first conclusion since the proof of the second conclusion follows similarly to that of Lemma 22(a). Write $h\left(\varepsilon_{i k}, \varepsilon_{j k}\right)=\mid$ $\left(\varepsilon_{i k}-\varepsilon_{j k}\right)-\left\langle\mathbf{e}_{k} \mathbf{x}_{i}^{\top}-\mathbf{e}_{k} \mathbf{x}_{j}^{\top}, \mathbf{A}-\mathbf{A}_{0}\right\rangle\left|-\left|\varepsilon_{i k}-\varepsilon_{j k}\right|\right.$,

$$
\mathcal{Z}_{M}=\sup _{\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1} \leq M}\left|\frac{1}{n(n-1)} \sum_{i \neq j} \sum_{k=1}^{m_{1}}\left\{h\left(\varepsilon_{i k}, \varepsilon_{j k}\right)-\mathbb{E}\left[h\left(\varepsilon_{i k}, \varepsilon_{j k}\right)\right]\right\}\right|
$$

For any $\gamma \geq 0$, Markov's inequality implies

$$
\mathbb{P}\left(\mathcal{Z}_{M}>t\right)=\mathbb{P}\left(e^{\gamma \mathcal{Z}_{M}}>e^{\gamma t}\right) \leq e^{-\gamma t} \mathbb{E} e^{\gamma \mathcal{Z}_{M}} .
$$

Let $M_{n}=\lfloor n / 2\rfloor$, the largest integer that is no larger than $n / 2$. Let $\left\{\sigma_{i k}\right\}, i=1, \ldots, n, k=$ $1, \ldots, m_{1}$ denote a Rademacher sequence independent of $\left\{\mathbf{x}_{i}, \boldsymbol{\varepsilon}_{i}\right\}, i=1, \ldots, n$. Similar to the proof of Lemma 22 we have

$$
\mathbb{E} e^{\gamma \mathcal{Z}_{M}} \leq \frac{1}{n!} \sum_{\pi} \mathbb{E} \exp \left(4 \gamma M\left\|\frac{1}{M_{n}} \sum_{i=1}^{M_{n}} \sum_{k=1}^{m_{1}} \sigma_{i k}\left(\mathbf{x}_{\pi(i)} \mathbf{e}_{k}^{\top}\left(m_{1}\right)-\mathbf{x}_{\pi\left(M_{n}+i\right)} \mathbf{e}_{k}^{\top}\left(m_{1}\right)\right)\right\|_{\mathrm{op}}\right) .
$$

To bound this term, using the $\varepsilon$-net discretization method in Lemma 22 and the subGaussion assumption on covariates, similarly we obtain

$$
\mathbb{E} e^{\gamma \mathcal{Z}_{M}} \leq 2 \cdot 9^{m_{1}+m_{2}} \exp \left(\frac{64 \gamma^{2} M^{2} \kappa_{0}^{2} C_{1}^{2}}{M_{n}}\right)
$$

for some constant $C_{1}$. Thus we have

$$
\begin{aligned}
\mathbb{P}\left(\mathcal{Z}_{M}>t\right) & \leq \inf _{\gamma \geq 0} e^{-\gamma t} \mathbb{E} e^{\gamma \mathcal{Z}_{M}} \\
& \leq \inf _{\gamma \geq 0} 2 \cdot 9^{m_{1}+m_{2}} \exp \left(C_{1}^{2} \kappa_{0}^{2} \gamma^{2} M^{2} / n-\gamma t\right) \\
& =2 \cdot 9^{m_{1}+m_{2}} \exp \left(-\frac{n t^{2}}{4 C_{1}^{2} \kappa_{0}^{2} M^{2}}\right) .
\end{aligned}
$$

Now we take $t=8 C_{1} \kappa_{0} \sqrt{\frac{m_{1}+m_{2}}{n}} M$, this implies $\mathcal{Z}_{M} \leq 8 C_{1} \kappa_{0} \sqrt{\frac{m_{1}+m_{2}}{n}} M$ with probability at least $1-2 \exp \left(-4\left(m_{1}+m_{2}\right)\right)$.
(b) The first conclusion follows from the definition of $\lambda^{*}$. Next we give the proof of the second conclusion. Conditional on $\left\{\mathbf{X}_{i}\right\}_{i=1}^{n}$. For any $\gamma \geq 0$, Markov's inequality implies

$$
\mathbb{P}\left(S_{n}>t \mid \mathbf{X}\right)=\mathbb{P}\left(e^{\gamma S_{n}}>e^{\gamma t} \mid \mathbf{X}\right) \leq e^{-\gamma t} \mathbb{E}_{\mathbb{P} \mid \mathbf{X}} e^{\gamma S_{n}} .
$$

By the $\varepsilon$-net discretization method and the sub-Gaussian assumption on covariates, we have

$$
\mathbb{E}_{\mathbb{P} \mid \mathbf{X}} e^{\gamma S_{n}} \leq 2 \cdot 9^{m_{1}+m_{2}} \exp \left(\frac{16 \gamma^{2}}{n} \max _{u \in \mathcal{N}, v \in \mathcal{M}} \frac{1}{n} \sum_{i=1}^{n}\left(u^{\top} \mathbf{x}_{i}\right)^{2}\right) .
$$

Assume $16 \max _{u \in \mathcal{N}} \frac{1}{n} \sum_{i=1}^{n}\left(u^{\top} \mathbf{x}_{i}\right)^{2} \leq \kappa_{0}^{2} C_{2}^{2}$ for some constant $C_{2}$, we have

$$
\begin{aligned}
\mathbb{P}\left(S_{n}>t \mid \mathbf{X}\right) & \leq \inf _{\gamma \geq 0} e^{-\gamma t} \mathbb{E}_{\mathbb{P} \mid \mathbf{X}} e^{\gamma S_{n}} \\
& \leq \inf _{\gamma \geq 0} 2 \cdot 9^{m_{1}+m_{2}} \exp \left(C_{2}^{2} \kappa_{0}^{2} \gamma^{2} / n-\gamma t\right) \\
& =2 \cdot 9^{m_{1}+m_{2}} \exp \left(-\frac{n t^{2}}{4 C_{2}^{2} \kappa_{0}^{2}}\right) .
\end{aligned}
$$

Now we take $t=4 C_{2} \kappa_{0} \sqrt{\frac{m_{1}+m_{2}}{n}}$, this implies $S_{n} \leq 4 C_{2} \kappa_{0} \sqrt{\frac{m_{1}+m_{2}}{n}}$ with probability at least $1-2 \exp \left(-\left(m_{1}+m_{2}\right)\right)$. Then if $m_{1}+m_{2}>\ln \left(2 / \alpha_{0}\right)$, we have $\lambda^{*}<4 C_{2} \kappa_{0} \sqrt{\frac{m_{1}+m_{2}}{n}}$.

Similar to the discussion of Lemma 22(b), we have that there exists constant $C_{2}$ and $c_{2}$ such that $16 \max _{u \in \mathcal{N}} \frac{1}{n} \sum_{i=1}^{n}\left(u^{\top} \mathbf{x}_{i}\right)^{2} \leq \kappa_{0}^{2} C_{2}^{2}$ with probability at least $1-9^{m_{1}+m_{2}} \exp \left(-c_{2} n\right)$. If $n>4\left(m_{1}+m_{2}\right) / c_{2}$, then we have $\lambda^{*}<4 C_{2} \kappa_{0} \sqrt{\frac{m_{1}+m_{2}}{n}}$ holds with probability at least $1-\exp \left(-\left(m_{1}+m_{2}\right)\right)$.
(c) We denote $\operatorname{Cov}\left(\mathbf{x}_{1}\right)$ by J. First we show that

$$
Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right) \geq \frac{b_{1}}{2}\left\|\mathbf{J}^{\frac{1}{2}}\left(\mathbf{A}-\mathbf{A}_{0}\right)^{\top}\right\|_{F}^{2} \wedge \frac{b_{1} b_{4}}{2}\left\|\mathbf{J}^{\frac{1}{2}}\left(\mathbf{A}-\mathbf{A}_{0}\right)^{\top}\right\|_{F}
$$

which implies that $\mathbf{A}_{0}$ is the minimizer of $Q(\mathbf{A})$. We define the maximal radius over which the criterion function can be minorated by a quadratic function
$r_{\mathbf{A}_{0}}=\sup _{r}\left\{r: Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right) \geq \frac{b_{1}}{2}\left\|\mathbf{J}^{1 / 2}\left(\mathbf{A}-\mathbf{A}_{0}\right)^{\top}\right\|_{F}^{2}\right.$ for all $\left.\left\|\mathbf{J}^{1 / 2}\left(\mathbf{A}-\mathbf{A}_{0}\right)^{\top}\right\|_{F} \leq r\right\}$.
We claim that $r_{\mathbf{A}_{0}} \geq b_{4}$. Note that the empirical loss of multivariate regression is a summation of $m_{1}$ empirical loss of linear regression, we have

$$
Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right) \geq b_{1}\left\|\mathbf{J}^{1 / 2}\left(\mathbf{A}-\mathbf{A}_{0}\right)^{\top}\right\|_{F}^{2}-\frac{b_{2}}{6} \sum_{k=1}^{m_{1}} \mathbb{E}_{\mathbf{X}}\left\{\left|\left(\mathbf{x}_{1}-\mathbf{x}_{2}\right)^{\top}\left(\mathbf{A}_{k}-\mathbf{A}_{0 k}\right)\right|^{3}\right\}
$$

Note that for all $\mathbf{A}-\mathbf{A}_{0}$, if

$$
\left\|\mathbf{J}^{1 / 2}\left(\mathbf{A}-\mathbf{A}_{0}\right)^{\top}\right\|_{F} \leq b_{4} \leq \frac{3 b_{1}}{2 \sqrt{2} b_{2}} \inf _{\mathbf{A}-\mathbf{A}_{0} \neq 0} \frac{\left(\sum_{k=1}^{m_{1}} \mathbb{E}\left|\left(\mathbf{x}_{i}-\mathbf{x}_{j}\right)^{\top}\left(\mathbf{A}_{k}-\mathbf{A}_{0 k}\right)\right|^{2}\right)^{3 / 2}}{\sum_{k=1}^{m_{1}} \mathbb{E}\left|\left(\mathbf{x}_{i}-\mathbf{x}_{j}\right)^{\top}\left(\mathbf{A}_{k}-\mathbf{A}_{0 k}\right)\right|^{3}},
$$

it follows that $\frac{b_{1}}{2}\left\|\mathbf{J}^{1 / 2}\left(\mathbf{A}-\mathbf{A}_{0}\right)^{\top}\right\|_{F}^{2} \geq \frac{b_{2}}{6} \sum_{k=1}^{m_{1}} \mathbb{E}_{\mathbf{X}}\left\{\left|\left(\mathbf{x}_{1}-\mathbf{x}_{2}\right)^{\top}\left(\mathbf{A}_{k}-\mathbf{A}_{0 k}\right)\right|^{3}\right\}$. Then we obtain

$$
Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right) \geq \frac{b_{1}}{2}\left\|\mathbf{J}^{1 / 2}\left(\mathbf{A}-\mathbf{A}_{0}\right)^{\top}\right\|_{F}^{2} \text { for all }\left\|\mathbf{J}^{1 / 2}\left(\mathbf{A}-\mathbf{A}_{0}\right)^{\top}\right\|_{F} \leq b_{4} .
$$

This implies $r_{\mathbf{A}_{0}} \geq b_{4}$. Similar to the proof of Lemma 22(c), we have

$$
Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right) \geq \frac{b_{1}}{2}\left\|\mathbf{J}^{\frac{1}{2}}\left(\mathbf{A}-\mathbf{A}_{0}\right)^{\top}\right\|_{F}^{2} \wedge \frac{b_{1} b_{4}}{2}\left\|\mathbf{J}^{\frac{1}{2}}\left(\mathbf{A}-\mathbf{A}_{0}\right)^{\top}\right\|_{F}
$$

for any $\mathbf{A}$. This implies that $\mathbf{A}_{0}$ is the minimizer of $Q(\mathbf{A})$, and

$$
Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right) \geq \frac{b_{1} b_{3}}{2}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{F}^{2} \wedge \frac{b_{1} b_{3} b_{4}}{2}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{F}
$$

for any $\mathbf{A}-\mathbf{A}_{0} \in \mathcal{C}$. Similar to the proof of Lemma $22(\mathrm{c})$, we can prove that

$$
Q(\widehat{\mathbf{A}})-Q\left(\mathbf{A}_{0}\right) \geq \frac{b_{1} b_{3}}{2}\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2}
$$

## Appendix E: Proof of Theorem 18

Hereafter we use $\widetilde{\mathbf{X}}$ and $\widetilde{\varepsilon}$ to denote $a \mathbf{X}$ and $a \varepsilon$ respectively, and let $\widetilde{\zeta}_{i j}=a_{i} \varepsilon_{i}-a_{j} \varepsilon_{j}$. With the following lemma and Theorem 1, we can get Theorem 18.

Lemma 25 Suppose that Assumptions 15-17 hold. Assume $\mathbf{A}_{0} \in \mathcal{B}_{q}\left(R_{q}\right) \cap \mathcal{S}$. Then we have following conclusions:
(a) Define for all $M>0$

$$
\mathcal{Z}_{M}:=\sup _{\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1} \leq M}\left|\left\{Q_{n}(\mathbf{A})-Q(\mathbf{A})\right\}-\left\{Q_{n}\left(\mathbf{A}_{0}\right)-Q\left(\mathbf{A}_{0}\right)\right\}\right| .
$$

Then we have for a constant $C_{0}>0$, with $\lambda_{\varepsilon}^{\prime}=8 \sqrt{2} C_{0} \sqrt{\frac{L \log \left(m_{1}+m_{2}\right)}{n m_{2}}}$ and $\tilde{\lambda}^{\prime}=$ $8 \eta \sqrt{\frac{\log \left(m_{1}+m_{2}\right)}{n}}$, such that

$$
\mathcal{Z}_{M} \leq \lambda_{\varepsilon}^{\prime} M+\tilde{\lambda}^{\prime}
$$

$\underset{\sim}{\text { with }} \underset{\sim}{\sim}$ probability at least $1-\exp \left(-4 \log \left(m_{1}+m_{2}\right)\right)$. Moreover, let $\lambda_{\varepsilon}=2 \lambda_{\varepsilon}^{\prime}$ and $\widetilde{\lambda}=2 \widetilde{\lambda^{\prime}}$, we have, for any $\mathbf{A} \in \mathcal{S}$

$$
\left|\left\{Q_{n}(\mathbf{A})-Q(\mathbf{A})\right\}-\left\{Q_{n}\left(\mathbf{A}_{0}\right)-Q\left(\mathbf{A}_{0}\right)\right\}\right| \leq \lambda_{\varepsilon}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1}+\widetilde{\lambda}
$$

with probability at least $1-\exp \left(-2 \log \left(m_{1}+m_{2}\right)\right)$.
(b) $\lambda^{*}>2\left\|\nabla Q_{n}\left(\mathbf{A}_{0}\right)\right\|_{\mathrm{op}}$ with probability $1-\alpha_{0}$ if $\varepsilon_{i}$ 's are symmetric. If $m_{1}+m_{2} \geq$ $\left(3 / \alpha_{0}\right)^{1 / 3}$, then $\lambda^{*} \leq \sqrt{\frac{5 L \log \left(m_{1}+m_{2}\right)}{n m_{2}}}$ with probability at least $1-\exp \left(-2 \log \left(m_{1}+m_{2}\right)\right)$.
(c) $\mathbf{A}_{0}$ is the minimizer of $Q(\mathbf{A})$. For all $\mathbf{A}, Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right) \geq \frac{1}{c_{1}^{2} \mu m_{1} m_{2}}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{F}^{2}$.

Proof [Proof of Lemma 25] (a) We only prove the first conclusion since the proof of the second conclusion follows similarly to that of Lemma $22(\mathrm{a})$. Due to the special structure of $\widetilde{\mathbf{X}}$, we adopt the bounded difference inequality to control the empirical process. Write $h\left(\widetilde{\varepsilon}_{i}, \widetilde{\varepsilon}_{j}\right)=\left|\left(\widetilde{\varepsilon}_{i}-\widetilde{\varepsilon}_{j}\right)-\left\langle\widetilde{\mathbf{X}}_{i}-\widetilde{\mathbf{X}}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle\right|-\left|\widetilde{\varepsilon}_{i}-\widetilde{\varepsilon}_{j}\right|$,

$$
\mathcal{Z}_{M}=\sup _{\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1} \leq M}\left|\frac{1}{n(n-1)} \sum_{i \neq j}\left\{h\left(\widetilde{\varepsilon}_{i}, \widetilde{\varepsilon}_{j}\right)-\mathbb{E}\left[h\left(\widetilde{\varepsilon}_{i}, \widetilde{\varepsilon}_{j}\right)\right]\right\}\right|
$$

Note that $\left|h\left(\widetilde{\varepsilon}_{i}, \widetilde{\varepsilon}_{j}\right)\right| \leq\left|\left\langle\widetilde{\mathbf{X}}_{i}-\widetilde{\mathbf{X}}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle\right| \leq 4 \eta$. Hence if we perturb one observation of the data set, the value of $\mathcal{Z}_{M}$ changes at $\frac{4 \eta}{n}$. By the bounded difference inequality, $\forall t>0$

$$
P\left(\mathcal{Z}_{M}-\mathbb{E}\left\{\mathcal{Z}_{M}\right\}>t\right) \leq \exp \left(-\frac{n t^{2}}{16 \eta^{2}}\right)
$$

Let $M_{n}=\lfloor n / 2\rfloor$, the largest integer that is no larger than $n / 2$. Let $\sigma_{i}, i=1, \ldots, n$ denote a Rademacher sequence independent of $\left\{\widetilde{\mathbf{X}}_{i}, \widetilde{\varepsilon}_{i}\right\}_{i=1}^{n}$, we have

$$
\begin{aligned}
\mathbb{E}\left[\mathcal{Z}_{M}\right] & =\mathbb{E}\left[\sup _{\|\mathbf{A}\|_{\infty} \leq \eta:\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1} \leq M} \frac{1}{n!}\left|\sum_{\pi} M_{n}^{-1} \sum_{i=1}^{M_{n}}\left\{h\left(\widetilde{\varepsilon}_{\pi(i)}, \widetilde{\varepsilon}_{\pi\left(M_{n}+i\right)}\right)-\mathbb{E} h\left(\widetilde{\varepsilon}_{\pi(i)}, \widetilde{\varepsilon}_{\pi\left(M_{n}+i\right)}\right)\right\}\right|\right] \\
& \leq 4 \mathbb{E}\left[\sup _{\|\mathbf{A}\|_{\infty} \leq \eta:\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1} \leq M}\left\|M_{n}^{-1} \sum_{i=1}^{M_{n}} \sigma_{i}\left(\widetilde{\mathbf{X}}_{\pi(i)}-\widetilde{\mathbf{X}}_{\pi\left(M_{n}+i\right)}\right)^{\top}\right\|_{\mathrm{op}}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{1}\right] \\
& \leq 4 M \mathbb{E}\left\|M_{n}^{-1} \sum_{i=1}^{M_{n}} \sigma_{i}\left(\widetilde{\mathbf{X}}_{\pi(i)}-\widetilde{\mathbf{X}}_{\pi\left(M_{n}+i\right)}\right)^{\top}\right\|_{\text {op }} \\
& \leq 8 \sqrt{2} M C_{0} \sqrt{\frac{L \log \left(m_{1}+m_{2}\right)}{n m_{2}}} \text { by triangle inequality and Lemma } 29 .
\end{aligned}
$$

Now we take $t=8 \eta \sqrt{\frac{\log \left(m_{1}+m_{2}\right)}{n}}$, this gives

$$
\mathcal{Z}_{M} \leq 8 \sqrt{2} M C_{0} \sqrt{\frac{L \log \left(m_{1}+m_{2}\right)}{n m_{2}}}+8 \eta \sqrt{\frac{\log \left(m_{1}+m_{2}\right)}{n}}
$$

with probability at least $1-\exp \left(-4 \log \left(m_{1}+m_{2}\right)\right)$.
(b) First we show that the symmetrized masks $a_{i} \mathbf{X}_{i}$ and noise $a_{i} \varepsilon_{i}$ are independent. Without loss of generality, we can calculate

$$
\begin{aligned}
& \mathbb{P}\left(a_{i} \mathbf{X}_{i}=\mathbf{e}_{1}\left(m_{1}\right) \mathbf{e}_{1}\left(m_{2}\right)^{\top}, a_{i} \varepsilon_{i}=\nu\right) \\
= & \mathbb{P}\left(a_{i}=1, \mathbf{X}_{i}=\mathbf{e}_{1}\left(m_{1}\right) \mathbf{e}_{1}\left(m_{2}\right)^{\top}, \varepsilon_{i}=\nu\right) \\
= & \mathbb{P}\left(a_{i}=1\right) \mathbb{P}\left(\mathbf{X}_{i}=\mathbf{e}_{1}\left(m_{1}\right) \mathbf{e}_{1}\left(m_{2}\right)^{\top}\right) \mathbb{P}\left(\varepsilon_{i}=\nu\right) \\
= & \mathbb{P}\left(a_{i}=1\right) \mathbb{P}\left(\mathbf{X}_{i}=\mathbf{e}_{1}\left(m_{1}\right) \mathbf{e}_{1}\left(m_{2}\right)^{\top}\right) \mathbb{P}\left(a_{i} \varepsilon_{i}=\nu\right) \\
= & \mathbb{P}\left(a_{i} \mathbf{X}_{i}=\mathbf{e}_{1}\left(m_{1}\right) \mathbf{e}_{1}\left(m_{2}\right)^{\top}\right) \mathbb{P}\left(a_{i} \varepsilon_{i}=\nu\right) .
\end{aligned}
$$

Then the first conclusion follows from the fact that $\left\{a_{i} \mathbf{X}_{i}\right\}_{i=1}^{n}$ is independent of $\left\{a_{i} \varepsilon_{i}\right\}_{i=1}^{n}$ and $\left\|\nabla Q_{n}\left(\mathbf{A}_{0}\right)\right\|_{\text {op }}$ has the same distribution as $S_{n}$ when $\varepsilon_{i}$ 's are symmetric.

Next we give the proof of the second conclusion. Conditional on $\left\{\widetilde{\mathbf{X}}_{i}\right\}_{i=1}^{n}$. For any $\gamma \geq 0$, Markov's inequality implies

$$
\mathbb{P}\left(S_{n}>t \mid \widetilde{\mathbf{X}}\right)=\mathbb{P}\left(e^{\gamma S_{n}}>e^{\gamma t} \mid \widetilde{\mathbf{X}}\right) \leq e^{-\gamma t} \mathbb{E}_{\mathbb{P} \mid \widetilde{\mathbf{X}}} e^{\gamma S_{n}}
$$

Recall that $S_{n}=\left\|\{n(n-1)\}^{-1} \sum_{i \neq j}\left(\widetilde{\mathbf{X}}_{j}-\widetilde{\mathbf{X}}_{i}\right) \operatorname{sign}\left(\varepsilon_{i}-\varepsilon_{j}\right)\right\|_{\text {op }}$, by the definition of operator norm, we have

$$
\begin{aligned}
& \mathbb{E}_{\mathbb{P} \mid \widetilde{\mathbf{X}}} \tilde{\mathrm{X}}^{\gamma_{n}} \\
= & \mathbb{E}_{\mathbb{P} \mid \widetilde{\mathbf{X}}} \exp \left(\gamma \sup _{\|u\|_{2} \leq 1,\|v\|_{2} \leq 1}\{n(n-1)\}^{-1} \sum_{i \neq j} u^{\top}\left(\widetilde{\mathbf{X}}_{i}-\widetilde{\mathbf{X}}_{j}\right) v \operatorname{sign}\left(\varepsilon_{i}-\varepsilon_{j}\right)\right) \\
= & \mathbb{E}_{\mathbb{P} \mid \widetilde{\mathbf{X}}} \exp \left(\gamma \sup _{\|u\|_{2} \leq 1,\|v\|_{2} \leq 1} \frac{1}{n!} \sum_{\pi} M_{n}^{-1} \sum_{i=1}^{M_{n}} u^{\top}\left(\widetilde{\mathbf{X}}_{\pi(i)}-\widetilde{\mathbf{X}}_{\pi\left(M_{n}+i\right)}\right) v \operatorname{sign}\left(\varepsilon_{\pi(i)}-\varepsilon_{\pi\left(M_{n}+i\right)}\right)\right) \\
\leq & \frac{1}{n!} \sum_{\pi} \mathbb{E}_{\mathbb{P} \mid \widetilde{\mathbf{X}}} \exp \left(\gamma \sup _{\|u\|_{2} \leq 1,\|v\|_{2} \leq 1} M_{n}^{-1} \sum_{i=1}^{M_{n}} u^{\top}\left(\widetilde{\mathbf{X}}_{\pi(i)}-\widetilde{\mathbf{X}}_{\pi\left(M_{n}+i\right)}\right) v \operatorname{sign}\left(\varepsilon_{\pi(i)}-\varepsilon_{\pi\left(M_{n}+i\right)}\right)\right) \\
= & \frac{1}{n!} \sum_{\pi} \mathbb{E}_{\mathbb{P} \mid \widetilde{\mathbf{X}}} \exp \left(\gamma\left\|M_{n}^{-1} \sum_{i=1}^{M_{n}}\left(\widetilde{\mathbf{X}}_{\pi(i)}-\widetilde{\mathbf{X}}_{\pi\left(M_{n}+i\right)}\right) \operatorname{sign}\left(\varepsilon_{\pi(i)}-\varepsilon_{\pi\left(M_{n}+i\right)}\right)\right\|_{\mathrm{op}}\right) .
\end{aligned}
$$

We denote $Z_{\pi}=\left\|M_{n}^{-1} \sum_{i=1}^{M_{n}}\left(\widetilde{\mathbf{X}}_{\pi(i)}-\widetilde{\mathbf{X}}_{\pi\left(M_{n}+i\right)}\right) \operatorname{sign}\left(\varepsilon_{\pi(i)}-\varepsilon_{\pi\left(M_{n}+i\right)}\right)\right\|_{\text {op }}$ for the simplicity of notation, then we have

$$
\mathbb{P}\left(S_{n}>t \mid \widetilde{\mathbf{X}}\right) \leq \inf _{\gamma \geq 0} e^{-\gamma t} \mathbb{E}_{\mathbb{P} \mid \widetilde{\mathbf{X}}} \tilde{e}^{\gamma S_{n}} \leq \inf _{\gamma \geq 0} e^{-\gamma t} \frac{1}{n!} \sum_{\pi} \mathbb{E}_{\mathbb{P} \mid \widetilde{\mathbf{X}}^{\gamma Z_{\pi}}}
$$

We do not use the $\varepsilon$-net argument to bound the operator norm term, instead we apply the Matrix Bernstein inequality (Theorem 6.1.1 in Tropp (2015)) to take advantage of the singleton design under the matrix completion setting. We state the Matrix Bernstein inequality in Lemma 26 for the sake of completeness. Now we calculate the quantity needed in Matrix Bernstein inequality. For all permutation $\pi$, we have

$$
\begin{aligned}
& \left\|M_{n}^{-1} \sum_{i=1}^{M_{n}}\left(\widetilde{\mathbf{X}}_{\pi(i)}-\widetilde{\mathbf{X}}_{\pi\left(M_{n}+i\right)}\right)\left(\widetilde{\mathbf{X}}_{\pi(i)}-\widetilde{\mathbf{X}}_{\pi\left(M_{n}+i\right)}\right)^{\top}\right\|_{\mathrm{op}} \\
= & \sup _{\|u\|_{2} \leq 1,\|v\|_{2} \leq 1} M_{n}^{-1} \sum_{i=1}^{M_{n}} u^{\top}\left(\widetilde{\mathbf{X}}_{\pi(i)}-\widetilde{\mathbf{X}}_{\pi\left(M_{n}+i\right)}\right)\left(\widetilde{\mathbf{X}}_{\pi(i)}-\widetilde{\mathbf{X}}_{\pi\left(M_{n}+i\right)}\right)^{\top} v
\end{aligned}
$$

By the special structure of $\widetilde{\mathbf{X}}$, we have $\pm u^{\top} \widetilde{\mathbf{X}}_{a} \widetilde{\mathbf{X}}_{b}^{\top} v \leq \frac{1}{2}\left(u^{\top} \widetilde{\mathbf{X}}_{a} \widetilde{\mathbf{X}}_{a}^{\top} u+v^{\top} \widetilde{\mathbf{X}}_{b} \widetilde{\mathbf{X}}_{b}^{\top} v\right)$, then for any $\|u\|_{2} \leq 1,\|v\|_{2} \leq 1$,

$$
\begin{aligned}
& M_{n}^{-1} \sum_{i=1}^{M_{n}} u^{\top}\left(\widetilde{\mathbf{X}}_{\pi(i)}-\widetilde{\mathbf{X}}_{\pi\left(M_{n}+i\right)}\right)\left(\widetilde{\mathbf{X}}_{\pi(i)}-\widetilde{\mathbf{X}}_{\pi\left(M_{n}+i\right)}\right)^{\top} v \\
\leq & \left(2 M_{n}\right)^{-1} \sum_{i=1}^{n}\left\{2 u^{\top} \widetilde{\mathbf{X}}_{i} \widetilde{\mathbf{X}}_{i}^{\top} v+u^{\top} \widetilde{\mathbf{X}}_{i} \widetilde{\mathbf{X}}_{i}^{\top} u+v^{\top} \widetilde{\mathbf{X}}_{i} \widetilde{\mathbf{X}}_{i}^{\top} v\right\} \\
= & 4 n^{-1}\left(\frac{u+v}{2}\right)^{\top} \sum_{i=1}^{n} \widetilde{\mathbf{X}}_{i} \widetilde{\mathbf{X}}_{i}^{\top}\left(\frac{u+v}{2}\right) .
\end{aligned}
$$

Note that $\left\|\frac{u+v}{2}\right\|_{2} \leq 1$, we obtain

$$
\left\|M_{n}^{-1} \sum_{i=1}^{M_{n}}\left(\widetilde{\mathbf{X}}_{\pi(i)}-\widetilde{\mathbf{X}}_{\pi\left(M_{n}+i\right)}\right)\left(\widetilde{\mathbf{X}}_{\pi(i)}-\widetilde{\mathbf{X}}_{\pi\left(M_{n}+i\right)}\right)^{\top}\right\|_{\mathrm{op}} \leq 4\left\|n^{-1} \sum_{i=1}^{n} \widetilde{\mathbf{X}}_{i} \widetilde{\mathbf{X}}_{i}^{\top}\right\|_{\mathrm{op}}
$$

Assume that $4\left\|n^{-1} \sum_{i=1}^{n} \widetilde{\mathbf{X}}_{i} \widetilde{\mathbf{X}}_{i}^{\top}\right\|_{\mathrm{op}} \leq \frac{8 L}{m_{2}}$, then we have

$$
\begin{aligned}
& \mathbb{E}_{\mathbb{P} \mid \widetilde{\mathbf{X}}} e^{\gamma Z_{\pi}} \\
= & \mathbb{E}_{\mathbb{P} \mid \widetilde{\mathbf{X}}} \int_{0}^{Z} \gamma e^{\gamma y} d y+1 \\
= & \int_{0}^{+\infty} \mathbb{P}\left(Z_{\pi}>y \mid \widetilde{\mathbf{X}}\right) \gamma e^{\gamma y} d y+1 \quad \text { by Fubini's theorem } \\
\leq & \gamma \int_{0}^{+\infty} \exp \left(-\frac{2 n^{2} y^{2}}{\frac{8 n L}{m_{2}}+n y}+\gamma y\right) d y+1 \quad \text { by Lemma } 26 \text { Matrix Bernstein inequality } \\
\leq & \gamma\left\{\int_{0}^{+\infty} \exp \left(-\frac{n m_{2} y^{2}}{8 L}+\gamma y\right) d y+\int_{0}^{+\infty} \exp (-n y+\gamma y) d y\right\}+1 \\
\leq & \gamma \sqrt{\frac{8 \pi L}{n m_{2}}} \exp \left(\frac{8 \gamma^{2} L}{4 n m_{2}}\right)+\int_{0}^{+\infty} \exp (-n y+\gamma y) d y+1
\end{aligned}
$$

Take $t=\sqrt{\frac{40 L \log \left(m_{1}+m_{2}\right)}{n m_{2}}}$ and $\gamma=n m_{2} t / 8 L$, if $n \gtrsim m_{2} \log \left(m_{1}+m_{2}\right)$, we obtain

$$
\begin{aligned}
& \mathbb{P}\left(S_{n}>t \mid \widetilde{\mathbf{X}}\right) \\
\leq & \inf _{\gamma \geq 0} e^{-\gamma t} \mathbb{E}_{\mathbb{P} \mid \widetilde{\mathbf{X}}} e^{\gamma Z_{\pi}} \\
\leq & e^{-5 \log \left(m_{1}+m_{2}\right)}\left\{\frac{n m_{2}}{8 L} \sqrt{\frac{40 L \log \left(m_{1}+m_{2}\right)}{n m_{2}}} \sqrt{\frac{8 \pi L}{n m_{2}}} \exp \left(\frac{5}{4} \log \left(m_{1}+m_{2}\right)\right)+\frac{1}{n-\gamma}+1\right\} \\
\leq & e^{-5 \log \left(m_{1}+m_{2}\right)}\left\{\sqrt{5 \pi \log \left(m_{1}+m_{2}\right)} \exp \left(\frac{5}{4} \log \left(m_{1}+m_{2}\right)\right)+2\right\} \\
\leq & 3 \exp \left(-3 \log \left(m_{1}+m_{2}\right)\right)
\end{aligned}
$$

thus if $m_{1}+m_{2} \geq\left(3 / \alpha_{0}\right)^{1 / 3}$, we conclude that $\lambda^{*} \leq \sqrt{\frac{40 L \log \left(m_{1}+m_{2}\right)}{n m_{2}}}$.
Finally we show that $4\left\|n^{-1} \sum_{i=1}^{n} \widetilde{\mathbf{X}}_{i} \widetilde{\mathbf{X}}_{i}^{\top}\right\|_{\mathrm{op}} \leq \frac{8 L}{m_{2}}$ with high probability. It's easy to see that $\left\|\mathbb{E} \widetilde{\mathbf{X}}_{i} \widetilde{\mathbf{X}}_{i}^{\top}\right\|_{\mathrm{op}} \leq \frac{L}{m_{2}}$, then by Matrix Bernstein inequality and $n \gtrsim m_{2} \log \left(m_{1}+m_{2}\right)$, we have

$$
\mathbb{P}\left(\left\|n^{-1} \sum_{i=1}^{n} \widetilde{\mathbf{X}}_{i} \widetilde{\mathbf{X}}_{i}^{\top}\right\|_{\mathrm{op}}>\frac{2 L}{m_{2}}\right) \leq\left(m_{1}+m_{2}\right) \exp \left(\frac{-\frac{L^{2}}{m_{2}^{2}}}{\frac{L}{n m_{2}}+\frac{1}{3 n} \frac{L}{m_{2}}}\right) \leq \exp \left(-2 \log \left(m_{1}+m_{2}\right)\right)
$$

Therefore, we conclude that $\lambda^{*} \leq \sqrt{\frac{5 L \log \left(m_{1}+m_{2}\right)}{n m_{2}}}$ with probability at least $1-\exp \left(-2 \log \left(m_{1}+m_{2}\right)\right)$.
(c) Similar to the proof of Lemma 22, we have

$$
\begin{aligned}
Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right)= & \mathbb{E}\left\langle\widetilde{\mathbf{X}}_{i}-\widetilde{\mathbf{X}}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle\left\{\mathrm{I}\left(\widetilde{\zeta}_{i j}<0\right)-1 / 2\right\} \\
& +\mathbb{E} \int_{0}^{\left\langle\widetilde{\mathbf{X}}_{i}-\widetilde{\mathbf{x}}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle}\left\{\mathrm{I}\left(\widetilde{\zeta}_{i j} \leq s\right)-\mathrm{I}\left(\widetilde{\zeta}_{i j} \leq 0\right)\right\} d s .
\end{aligned}
$$

For the first part,

$$
\begin{aligned}
& \mathbb{E}\left\langle\widetilde{\mathbf{X}}_{i}-\widetilde{\mathbf{X}}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle\left\{\mathrm{I}\left(\widetilde{\zeta}_{i j}<0\right)-1 / 2\right\} \\
= & \frac{1}{2} \mathbb{E}\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle\left\{\mathrm{I}\left(\zeta_{i j}^{-}<0\right)-1 / 2\right\}+\frac{1}{2} \mathbb{E}\left\langle\mathbf{X}_{i}+\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle\left\{\mathrm{I}\left(\zeta_{i j}^{+}<0\right)-1 / 2\right\} \\
= & 0 .
\end{aligned}
$$

For the second part,

$$
\begin{aligned}
& \mathbb{E} \int_{0}^{\left\langle\widetilde{\mathbf{X}}_{i}-\widetilde{\mathbf{x}}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle}\left\{\mathrm{I}\left(\widetilde{\zeta}_{i j} \leq s\right)-\mathrm{I}\left(\widetilde{\zeta}_{i j} \leq 0\right)\right\} d s \\
= & \frac{1}{2} \mathbb{E} \int_{0}^{\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle}\left\{\mathrm{I}\left(\zeta_{i j}^{-} \leq s\right)-\mathrm{I}\left(\zeta_{i j}^{-} \leq 0\right)\right\} d s \\
& +\frac{1}{2} \mathbb{E} \int_{0}^{\left\langle\mathbf{X}_{i}+\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle}\left\{\mathrm{I}\left(\zeta_{i j}^{+} \leq s\right)-\mathrm{I}\left(\zeta_{i j}^{+} \leq 0\right)\right\} d s .
\end{aligned}
$$

Denote the distribution function of $\zeta_{i j}^{-}$by $F^{*}(\cdot)$. By the independence of $\mathbf{X}$ and $\varepsilon$,

$$
\begin{aligned}
& \frac{1}{2} \mathbb{E} \int_{0}^{\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle}\left\{\mathrm{I}\left(\zeta_{i j}^{-} \leq s\right)-\mathrm{I}\left(\zeta_{i j}^{-} \leq 0\right)\right\} d s \\
= & \frac{1}{2} \mathbb{E}_{\mathbf{X}_{i}, \mathbf{X}_{j}} \int_{0}^{\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle}\left\{F^{*}(s)-F^{*}(0)\right\} d s \\
= & \frac{1}{2} \mathbb{E}_{\mathbf{X}_{i}, \mathbf{X}_{j}} \int_{0}^{\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle} f^{-}\left(\xi_{i j}\right) s d s \\
\geq & \frac{1}{4 c_{1}^{2}} \mathbb{E}_{\mathbf{X}_{i}, \mathbf{X}_{j}}\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle^{2} .
\end{aligned}
$$

For the same reason, we have

$$
\frac{1}{2} \mathbb{E} \int_{0}^{\left\langle\mathbf{X}_{i}+\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle}\left\{\mathrm{I}\left(\zeta_{i j}^{+} \leq s\right)-\mathrm{I}\left(\zeta_{i j}^{+} \leq 0\right)\right\} d s \geq \frac{1}{4 c_{1}^{2}} \mathbb{E}_{\mathbf{X}_{i}, \mathbf{X}_{j}}\left\langle\mathbf{X}_{i}+\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle^{2} .
$$

From the above discussion, we conclude that

$$
\begin{aligned}
Q(\mathbf{A})-Q\left(\mathbf{A}_{0}\right) & \geq \frac{1}{4 c_{1}^{2}} \mathbb{E}_{\mathbf{X}_{i}, \mathbf{X}_{j}}\left\langle\mathbf{X}_{i}-\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle^{2}+\frac{1}{4 c_{1}^{2}} \mathbb{E}_{\mathbf{X}_{i}, \mathbf{X}_{j}}\left\langle\mathbf{X}_{i}+\mathbf{X}_{j}, \mathbf{A}-\mathbf{A}_{0}\right\rangle^{2} . \\
& =\frac{1}{c_{1}^{2}} \mathbb{E}_{\mathbf{X}_{i}}\left\langle\mathbf{X}_{i}, \mathbf{A}-\mathbf{A}_{0}\right\rangle^{2} \\
& \geq \frac{1}{c_{1}^{2} \mu m_{1} m_{2}}\left\|\mathbf{A}-\mathbf{A}_{0}\right\|_{F}^{2},
\end{aligned}
$$

which also implies that $\mathbf{A}_{0}$ is the minimizer of $Q(\mathbf{A})$.

## Appendix F: Collection of Lemmas

In this part, we give some lemmas used in the previous proofs.
Lemma 26 (Matrix Bernstein, Theorem 6.1.1 in Tropp (2015)) Consider a finite sequence $\left\{\mathbf{Z}_{i}\right\}_{i=1}^{n}$ of independent, random matrices with common dimension $m_{1} \times m_{2}$. Assume that $\mathbb{E} \mathbf{Z}_{i}=\mathbf{0}$ and $\left\|\mathbf{Z}_{i}\right\|_{\mathrm{op}} \leq H$ for each index $k$, introduce the random matrix

$$
\mathbf{S}=\sum_{k} \mathbf{Z}_{i} .
$$

Let $v(\mathbf{S})$ be the matrix variance statistic of the sum:

$$
\begin{aligned}
v(\mathbf{S}) & =\max \left\{\left\|\mathbb{E}\left(\mathbf{S S}^{\top}\right)\right\|_{\mathrm{op}},\left\|\mathbb{E}\left(\mathbf{S}^{\top} \mathbf{S}\right)\right\|_{\mathrm{op}}\right\} \\
& =\max \left\{\left\|\sum_{k} \mathbb{E}\left(\mathbf{Z}_{i} \mathbf{Z}_{i}^{\top}\right)\right\|_{\mathrm{op}},\left\|\sum_{k} \mathbb{E}\left(\mathbf{Z}_{i}^{\top} \mathbf{Z}_{i}\right)\right\|_{\mathrm{op}}\right\} .
\end{aligned}
$$

Then for all $t \geq 0$,

$$
\mathbb{P}\left\{\|\mathbf{S}\|_{\mathrm{op}} \geq t\right\} \leq\left(m_{1}+m_{2}\right) \exp \left(\frac{-t^{2} / 2}{v(\mathbf{S})+H t / 3}\right)
$$

Matrix Bernstein requires the constant bound on $\left\|\mathbf{Z}_{i}\right\|_{\text {op }}$, it's possible to replace this requirement by the bound on the weaker $\psi_{p}$-norms of $\left\|\mathbf{Z}_{i}\right\|_{\text {op }}$, we have following result which are very useful in matrix completion problem. Let $\mathbf{Z}_{1}, \ldots, \mathbf{Z}_{n}$ be independent random matrices with dimensions $m_{1} \times m_{2}$. Define

$$
\sigma_{\mathbf{Z}}^{2}=\max \left\{\left\|\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left(\mathbf{Z}_{i} \mathbf{Z}_{i}^{\top}\right)\right\|_{\mathrm{op}},\left\|\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left(\mathbf{Z}_{i}^{\top} \mathbf{Z}_{i}\right)\right\|_{\mathrm{op}}\right\}
$$

Lemma 27 (Proposition 2 in Koltchinskii et al. (2011)) Let $\mathbf{Z}, \mathbf{Z}_{1}, \ldots, \mathbf{Z}_{n}$ be i.i.d. random matrices with dimensions $p \times q$ that satisfy $\mathbb{E}(\mathbf{Z})=0$. Suppose that $\left\|\|\mathbf{Z}\|_{\text {op }}\right\|_{\psi_{p}}<\infty$ for some $p \geq 1$. Then there exists a constant $C>0$ such that, for all $t>0$, with probability at least $1-e^{-t}$

$$
\begin{aligned}
\left\|\frac{\mathbf{Z}_{1}+\cdots+\mathbf{Z}_{n}}{n}\right\|_{\mathrm{op}} \leq C \max & \left\{\sigma \mathbf{Z} \sqrt{\frac{t+\log \left(m_{1}+m_{2}\right)}{n}}\right. \\
& \left.\left\|\|\mathbf{Z}\|_{\mathrm{op}}\right\|_{\psi_{p}}\left(\log \frac{\| \| \mathbf{Z}\left\|_{\mathrm{op}}\right\|_{\psi_{p}}}{\sigma_{\mathbf{Z}}}\right)^{1 / p} \frac{t+\log \left(m_{1}+m_{2}\right)}{n}\right\} .
\end{aligned}
$$

Lemma 28 Let $\left\{\mathbf{X}_{i}\right\}_{i=1}^{n}$ be i.i.d. with distribution $\Pi$ on $\mathcal{X}$ which satisfies Assumptions 15 and 16. Assume that $\left(\sigma_{i}\right)_{i=1}^{n}$ are i.i.d. Rademacher random variables independent of $\left\{\mathbf{X}_{i}\right\}_{i=1}^{n}$. Then, there exists an absolute constant $C^{*}>0$ that depends only on $K$ and such that, for all $t>0$ with probability at least $1-e^{-t}$ we have

$$
\left\|\frac{1}{n} \sum_{i=1}^{n} \sigma_{i} \mathbf{X}_{i}\right\|_{\mathrm{op}} \leq C^{*} \max \left\{\sqrt{\frac{L\left(t+\log \left(m_{1}+m_{2}\right)\right)}{n m_{2}}}, \sqrt{\log m_{2}} \frac{\left(t+\log \left(m_{1}+m_{2}\right)\right)}{n}\right\}
$$

Proof We apply Lemma 27 to $\mathbf{Z}_{i}=\sigma_{i} \mathbf{X}_{i}$. We first estimate $\sigma_{\mathbf{Z}}$ and $\left\|\|\mathbf{Z}\|_{\text {op }}\right\|_{\psi_{2}}$. Note that $\mathbf{Z}_{i}$ is a zero-mean random matrix which satisfies $\left\|\mathbf{Z}_{i}\right\|_{\text {op }} \leq\left|\sigma_{i}\right|$, hence $\left\|\left\|\mathbf{Z}_{i}\right\|_{\text {op }}\right\|_{\psi_{2}} \leq\left\|\sigma_{i}\right\|_{\psi_{2}}$, together with the fact that the $\psi_{2}$ norm of a Rademacher random variable $\sigma_{i}$ is equal to $\left\|\sigma_{i}\right\|_{\psi_{2}}=\sqrt{1 / \log 2}$, we have $\left\|\|\mathbf{Z}\|_{\text {op }}\right\|_{\psi_{2}} \leq \sqrt{1 / \log 2}$. Then we compute $\mathbb{E}\left(\mathbf{Z}_{i} \mathbf{Z}_{i}^{\top}\right)=R$ and $\mathbb{E}\left(\mathbf{Z}_{i}^{\top} \mathbf{Z}_{i}\right)=C$, where $C$ and $R$ is the diagonal matrix with $C_{k}$ and $R_{j}$ on the diagonal respectively. This and the fact that the $\mathbf{Z}_{i}$ are i.i.d. imply that $\sigma_{\mathbf{Z}}^{2}=\max _{i, j}\left(C_{i}, R_{j}\right) \leq$ $L / m_{2}$. Note that $\max _{i, j}\left(C_{i}, R_{j}\right) \geq 1 / q$ together with the concavity of the logarithm imply that $\sqrt{\log \left(\left\|\|\mathbf{Z}\|_{\text {op }}\right\|_{\psi_{p}} / \sigma_{\mathbf{Z}}\right)} \leq \sqrt{\log m_{2}}$. Hence the statement of Lemma 28 follows.
We choose the parameter $t$ to be $t=\log \left(m_{1}+m_{2}\right)$. Next we use Lemma 28 to bound the expectation of the largest singular value of the sum of masks.

Lemma 29 Let $\left\{\mathbf{X}_{i}\right\}_{i=1}^{n}$ be i.i.d. with distribution $\Pi$ on $\mathbf{X}$ which satisfies Assumptions 15 and 16. Assume that $\left\{\sigma_{i}\right\}_{i=1}^{n}$ are i.i.d. Rademacher random variables independent of $\left\{\mathbf{X}_{i}\right\}_{i=1}^{n}$. Then, for $n>2 m_{2} \log ^{2}\left(m_{1}+m_{2}\right) / L$, there exists an absolute constant $C_{0}>0$ such that

$$
\mathbb{E}\left\|\frac{1}{n} \sum_{i=1}^{n} \sigma_{i} \mathbf{X}_{i}\right\|_{\mathrm{op}} \leq C_{0} \sqrt{\frac{L \log \left(m_{1}+m_{2}\right)}{n m_{2}}} .
$$

Proof This lemma is obtained by integrating the tail probability of Lemma 28.

## Appendix G: Algorithms and Complexity Analysis

G1: Proximal Gradient Algorithm for the Rank Matrix Lasso
In this subsection, we propose an accelerated proximal gradient (APG) algorithm to solve the minimization (4). Specifically, to minimize a penalized loss function, i.e.,

$$
\begin{equation*}
\min _{\mathbf{A}}\left\{Q_{n}(\mathbf{A})+\lambda\|\mathbf{A}\|_{1}\right\}, \tag{G.1}
\end{equation*}
$$

we employ the quadratic function

$$
\begin{equation*}
Q_{\text {Major }}\left(\mathbf{A} ; \mathbf{A}^{(l)}\right)=Q_{n}\left(\mathbf{A}^{(l)}\right)+\left\langle\nabla Q_{n}\left(\mathbf{A}^{(l)}\right), \mathbf{A}-\mathbf{A}^{(l)}\right\rangle+(L / 2)\left\|\mathbf{A}-\mathbf{A}^{(l)}\right\|_{F}^{2} \tag{G.2}
\end{equation*}
$$

to locally "majorize" $Q_{n}(\mathbf{A})$ for the $t$-th iteration (Fan et al., 2018), where $L$ is a constant such that in a neighborhood of $\mathbf{A}^{(l)}$ we have $Q_{n}(\mathbf{A}) \leq Q_{\text {Major }}\left(\mathbf{A} ; \mathbf{A}^{(l)}\right)$ and the equality can be attained at $\mathbf{A}^{(l)}$. Then, we solve

$$
\begin{equation*}
\min _{\mathbf{A}} L_{\text {Major }}\left(\mathbf{A} ; \mathbf{A}^{(l)}\right)=\min _{\mathbf{A}}\left\{Q_{\text {Major }}\left(\mathbf{A} ; \mathbf{A}^{(l)}\right)+\lambda\|\mathbf{A}\|_{1}\right\} \tag{G.3}
\end{equation*}
$$

and set the solution as $\mathbf{A}^{(l+1)}$, which gives

$$
Q_{n}\left(A^{(l+1)}\right)+\lambda\left\|A^{(l+1)}\right\|_{1} \leq L_{\text {Major }}\left(A^{(l+1)} ; A^{(l)}\right) \leq L_{\text {Major }}\left(A^{(l)} ; A^{(l)}\right)=Q_{n}\left(A^{(l)}\right)+\lambda\left\|A^{(l)}\right\|_{1} .
$$

While typically (G.1) does not have a closed-form solution, the minimizer in (G.3) can be expressed using the singular value soft-thresholding operator (see e.g. Toh and Yun
(2010)): for any given $\mathbf{M} \in \mathbb{R}^{m_{1} \times m_{2}}, \operatorname{Soft}(\mathbf{M} ; \tau)=\mathbf{U}_{\mathbf{M}}^{\top} \operatorname{diag}\left\{\left(\sigma_{i}(\mathbf{M})-\tau\right)_{+}\right\} \mathbf{V}_{\mathbf{M}}$. Here $\mathbf{M}=\mathbf{U}_{\mathbf{M}}^{\top} \operatorname{diag}\left\{\left(\sigma_{i}(\mathbf{M})\right)\right\} \mathbf{V}_{\mathbf{M}}$ is the singular value decomposition of $\mathbf{M}$. Detailed description and explicit mathematical expressions are provided in Algorithm 1. In our simulations, the tol we picked is $10^{-4}$ for all the trials, and the maximal iteration $T=100$. We take $L^{(0)}=10^{-4} L_{\max }$, and set $L_{\max }=3 \times 10^{p}$ which is empirically found good enough for convergence in most scenarios.

```
Algorithm 1: Accelerated proximal gradient algorithm for the rank matrix lasso
    Input: Observed data \(\left(\mathbf{y}_{i}, \mathbf{X}_{i}\right)\), for \(i=1, \cdots, n\); tuning parameter \(\lambda^{\star}\) by (8); floor
                curvature \(L^{(0)}\); ceiling curvature \(L_{\text {max }}\); updating rate \(\eta ; t^{(0)}=t^{(1)}=1\);
                convergence tolerance tol and maximal iteration \(T\); initial estimate
                    \(\mathbf{A}^{(0)}=\mathbf{A}^{(1)}=\mathbf{0}\).
    Output: Estimator \(\widehat{\mathbf{A}}\).
    Set \(\mathbf{B}^{(l)}=\mathbf{A}^{(l)}+\frac{t^{(l-1)}-1}{t^{(l)}}\left(\mathbf{A}^{(l)}-\mathbf{A}^{(l-1)}\right)\).
    Calculate the sub-gradient \(\mathbf{G}^{(l)}=\left.\nabla Q_{n}(\mathbf{A})\right|_{\mathbf{A}=\mathbf{B}^{(l)}}\).
    Set \(L=\min \left\{\eta L^{(l-1)}, L_{\max }\right\}\).
    while \(L<L_{\text {max }}\) do
        Compute \(\mathbf{S}=\operatorname{Soft}\left(\mathbf{B}^{(l)}-L^{-1} \mathbf{G}^{(l)} ; L^{-1} \lambda^{\star}\right)\);
        Calculate \(Q_{n}(\mathbf{S})\) and \(Q_{\text {Major }}\left(\mathbf{S} ; \mathbf{B}^{(l)}\right)\);
        if \(Q_{n}(\mathbf{S}) \leq Q_{\text {Major }}\left(\mathbf{S} ; \mathbf{B}^{(l)}\right)\) then
            Set \(\mathbf{A}^{(l+1)}=\mathbf{S}, L^{(l)}=L\)
            break
        else
            Set \(L=\min \left\{L / \eta, L_{\text {max }}\right\}\).
12 Compute \(t^{(l+1)}=\frac{1+\sqrt{1+4\left(t^{(l)}\right)^{2}}}{2}\).
13 Repeat above steps until the stop criterion is meet:
    \(\left\|\mathbf{A}^{(l+1)}-\mathbf{A}^{(l)}\right\|_{F} /\left\|\mathbf{A}^{(l)}\right\|_{F} \leq\) tol or the maximal number of iteration \(T\) is hit. Set
    \(\widehat{\mathbf{A}}=\mathbf{A}^{(l+1)}\).
```

Remark 30 The suggested APG is kind of a "first-order" algorithm, while some secondorder algorithms, such as quasi-Newton algorithms, are available for tackling nuclear norm penalization problems. In Appendix H2, we implement a state-of-the-art quasi-Newton algorithm proposed by Becker et al. (2019) for solving the rank-based optimization and make a comparison with Algorithm 1. An interesting trade-off is that, while the second-order method has certain advantage on the total steps for convergence, it generally requires more time in each step since first-order methods have to be implied to solve the subproblems. Our empirical results reveal that the proposed $A P G$ algorithm is at least comparable with the quasi-Newton algorithm in terms of average computation time.

Remark 31 While "smooth + nonsmooth" optimization has been extensively studied (Nesterov, 2013; Lee et al., 2014; Chambolle and Dossal, 2015), the global convergence rates of the APG method for "nonsmooth + nonsmooth" optimization has not been fully understood (Bian and Wu, 2021). Current theoretical progress on this issue typically involves
some modification on the procedure. For example, Yu et al. (2010) uses a "tightest pseudo quadratic fit" in the proximal step, which gives global convergence in objective function value. However, this would render the subproblem hard to solve, especially for the nuclear norm penalization. Another popular direction is to locally approximate the the loss function by smoothing methods (Nesterov, 2005; Zhang and Chen, 2009; Bian and Chen, 2020; Bian and $W u$, 2021), but a careful tuning procedure is generally required. It is still not clear whether these solutions with certain guarantee of global convergence can be extended to our objective function and tuning scheme, but this certainly warrants future research.

## G2: Complexity Analysis

We present an analysis of the time complexity of the entire algorithm including both parts of parameter selection and optimizations. We take the trace regression as an example for careful investigation, but other models such as multivariate regression can be studied analogously. Our standpoint is the primitive form of the algorithm without taking special structural blessing (such as sparse matrices or factor matrices) into consideration. For trace regression models, the complexity of Algorithm 1 can be decomposed into the following steps:

Pivotal tuning. Computing the summation of the gradient is of order $O\left(n m_{1} m_{2}\right)$. Uses SVD for obtaining the operator norm of $\nabla Q_{n}(\mathbf{A})$ requires $O\left(m_{1} m_{2} \min \left\{m_{1}, m_{2}\right\}\right)$. It is usually redundant to perform a full SVD. There are other algorithms that are less expensive, like power iteration or Lanczos bidiagonalization (Baglama and Reichel, 2005; Larsen, 2004). Hence, the total complexity with $B$ rounds of simulations is about $O\left\{B\left(n m_{1} m_{2}+m_{1}^{2} m_{2}\right)\right\}$ (assume $m_{1} \leq m_{2}$ ). In most problems, $n$ is typically required to be larger than $m_{1}$ and $m_{2}$ to achieve successful recovery. Therefore, the complexity for Pivotal tuning is basically $O\left(B n m_{1} m_{2}\right)$.

APG optimization The following steps contribute most significantly to the computation costs: (1) Subgradient calculation in Step 2, which by our previous argument, takes $O\left(n m_{1} m_{2}\right)$ operations for the problems of interest. (2) SVD in Step 5. In medium scale problems full SVD can be directly applied, which takes $O\left(m_{1} m_{2} \min \left\{m_{1}, m_{2}\right\}\right)$ operations. Again, utilizing other algorithms for structured problems could overcome this barrier and accelerate the program significantly. (3) Calculation of the objective value and the majorized value in Step 6. The former involves a sorting of the residuals in $O(n \log n)$ operations and summation of $n$ linear products. The latter demands a calculation of $Q_{n}\left(\mathbf{B}^{(l)}\right)$ as well as two inner products for the first-order and quadratic approximation terms (see (G.2)). Thus the dominant part takes up to $O\left(n m_{1} m_{2}\right)$ operations.

From the above analysis we can conclude that APG proceeds with the complexity of $O\left(n m_{1} m_{2}\right)$ each iteration, which aligns well with the Pivotal tuning step. Our simulation provides evidence that the time for parameter selection and for optimization remains in a similar magnitude among various models and settings. Empirically, the pivotal tuning usually takes less time than the APG program. These results justify the superiority of pivotal tuning compared with some conventional methods like cross-validation in matrix settings.

In Appendix H1, we also discuss how to make our algorithm scalable to large-scale structured matrix recovery problems and verify its effectiveness via a simulation study.

## Appendix H: Additional Simulations Results

## H1: Adaptivity to Large-Scale Problems: Large-Scale Matrix Completion

Real-life application typically involves computation of large matrices, especially for matrix completion problem. By taking advantage of the structure of the problem(like sparsity in matrix completion scenarios) as well as state-of-art algorithms for handling operator norm and SVD computation(Larsen, 1998; Baglama and Reichel, 2005; Larsen, 2004), our program can be adapted to large-scale structures and demonstrate the computational efficiency gain of pivotal tuning. For example, large-scale matrix completion problems are endowed with a perfect sparse structure. Then instead of the $O\left(n m_{1} m_{2}\right)$ costs we only need $O(n)$ operations for calculating the sub-gradients and objective values(with $O(n \log n)$ rounds of sorting), which is undoubtedly a striking improvement over the direct summation. Furthermore, instead of performing full SVD each time, a more appealing trick is to gradually increase and explore the rank as the program progresses(Toh and Yun, 2010). With this strategy we typically only need to focus on several largest eigen-pairs in each iteraction, for which many fast algorithms exist. In our large-scale implementation we use the PROPACK package by Larsen (2004) to achieve this convenience, rendering our algorithm perfectly scalable to large matrix computation problems $\left(m_{1}, m_{2}\right.$ as high as $\left.5 \times 10^{4}\right)$. As a concluding remark, some other improvements are also possible in large-scale extensions. For example, it is appealing to take advantage of robust distributed programs and parallelization algorithms. See for example Chen et al. (2020).

Here we report some numerical results on large-scale matrix completion problems. For operator norm and SVD calculation we make use of the lansvd function in PROPACK. We consider $m=5 \times 10^{3}, 1 \times 10^{4}, 5 \times 10^{4}$ and $r=10$. The ground truth $\mathbf{A}_{0}$ and the noise are generated similarly to our small-scale matrix completion setting. Each time we sample $N=6 \cdot d f$ observations, where the degree of freedom of a matrix $d f=r\left(m_{1}+\right.$ $\left.m_{2}-r\right)=r(2 m-r)$. Table S. 3 summarised the results. The results demonstrate both the computational efficiency and the statistical accuracy of the proposed algorithm.

| Noise | Dimension | $\left\\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\\|_{F}^{2}$ | Rank | Tuning(s) | Solving(s) | Iteration |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| Gaussian | $5 \times 10^{3}$ | $0.180(8.7 \mathrm{e}-4)$ | $10.0(0)$ | 4.77 | 13.8 | 39.9 |
|  | $1 \times 10^{4}$ | $0.180(6.9 \mathrm{e}-4)$ | $10.0(0)$ | 9.85 | 26.8 | 36.4 |
|  | $5 \times 10^{4}$ | $0.179(3.4 \mathrm{e}-4)$ | $10.0(0)$ | 57.4 | 225 | 32.0 |
| Cauchy | $5 \times 10^{3}$ | $0.004(6.6 \mathrm{e}-5)$ | $10.0(0)$ | 4.74 | 34.5 | 84.0 |
|  | $1 \times 10^{4}$ | $0.004(7.5 \mathrm{e}-5)$ | $10.0(0)$ | 9.86 | 70.5 | 84.4 |
|  | $5 \times 10^{4}$ | $0.004(2.0 \mathrm{e}-5)$ | $10.0(0)$ | 58.7 | 723 | 87.4 |

Table S.3: Simulation for large-scale matrix completion

## H2: Comparison of First-Order and Second-Order Algorithms

Our main trials are based on a first-order algorithm. The population version of the proposed loss function has higher order curvature which is coarsely represented by a constant diagonal matrix in first-order methods. We now frame a quasi-newton algorithm for minimizing the
rank-based objective, attempting to borrow information from second-order geometry of the formulation.

```
Algorithm 2: Proximal quasi-Newton algorithm for the rank matrix lasso
    Input: Observed data \(\left(\mathbf{y}_{i}, \mathbf{X}_{i}\right)\), for \(i=1, \cdots, n\); pivotal tuning parameter \(\alpha\);
            updating rate \(\eta\); convergence tolerance tol and maximal iteration \(T\); initial
            estimate \(\mathbf{A}^{(0)}\); tuning parameter \(L\).
    Output: Estimator \(\widehat{\mathbf{A}}\).
    /* Pivotal tuning
                                    */
1 Calculate \(\lambda^{\star}\) using the pivotal tuning scheme with confidence level \(\alpha\).
    /* Start optimization */
    Repeat
    3 Choose \(\mathbb{B}^{(l)}\) to be a Hessian approximation at \(\mathbf{A}^{(l)}\).
    4 Solve the subproblem for a searching direction using a first-order forward-backward
    splitting algorithm:
\[
\Delta \mathbf{A}^{(l)} \leftarrow \arg \min _{\mathbf{D}} Q_{n}\left(\mathbf{A}^{(l)}\right)+\left\langle\nabla Q_{n}\left(\mathbf{A}^{(l)}\right), \mathbf{D}\right\rangle+\left(\mathbf{D}^{V}\right)^{\top} \mathbb{B}^{(l)} \mathbf{D}^{V}+\lambda^{\star}\left\|\mathbf{A}^{(l)}+\mathbf{D}\right\|_{1} .
\]
5 Select \(t^{(l)}\) with a backtracking line search;
6 Update: \(\mathbf{A}^{(l+1)}=\mathbf{A}^{(l)}+t^{(l)} \Delta \mathbf{A}^{(l)}\).
\(\mathbf{7}\) until: \(\left\|\mathbf{A}^{(l+1)}-\mathbf{A}^{(l)}\right\|_{F} /\left\|\mathbf{A}^{(l)}\right\|_{F} \leq\) tol or the maximal number of iteration \(T\) is hit.
\(8 \operatorname{Set} \widehat{\mathbf{A}}=\mathbf{A}^{(l+1)}\).
```

We consider several details regarding the implementation of the algorithm:

- Steps to update the Hessian approximation matrix(Step 3 in Algorithm 2). Since we are targeting problems involving large matrices, it is a necessity to avoid saving or loading dense Hessian matrices. To this end, we pursue using limited memories Hessian approximation schemes. Two simple yet admirable options in the literature are SR1 update(Becker and Fadili, 2012; Becker et al., 2019) and BFGS update(Lee et al., 2014), which provide rank-1 and rank-2 updates respectively. Concretely speaking, let

$$
\mathbf{d} \mathbf{S}^{(l)}=\left\{\mathbf{A}^{(l)}-\mathbf{A}^{(l-1)}\right\}^{V}, \mathbf{d} \mathbf{G}^{(l)}=\left\{\nabla Q_{n}\left(\mathbf{A}^{(l)}\right)-\nabla Q_{n}\left(\mathbf{A}^{(l-1)}\right)\right\}^{V}
$$

The 1-step SR1 update gives:

$$
\mathbb{B}^{(l)}=\mathbb{B}_{0}^{(l)}+\mathbf{u}^{(l)} \mathbf{u}^{(l) \top}
$$

where

$$
\mathbf{u}^{(l)}=\frac{\mathbf{v}^{(l)} \mathbf{v}^{(l) \top}}{\mathbf{v}^{(l)^{\top}} \mathbf{d} \mathbf{S}^{(l)}}, \mathbf{v}^{(l)}=\mathbf{d} \mathbf{G}^{(l)}-\mathbb{B}^{(l)} \mathbf{d} \mathbf{S}^{(l)}
$$

The 1-step BFGS gives:

$$
\mathbb{B}^{(l)}=\mathbb{B}_{0}^{(l)}+\mathbf{u}_{\mathbf{1}}^{(l)} \mathbf{u}_{\mathbf{1}}^{(l) \top}-\mathbf{u}_{\mathbf{2}}^{(l)} \mathbf{u}_{\mathbf{2}}^{(l) \top}
$$

where

$$
\mathbf{u}_{1}^{(l)}=\frac{\mathbf{d G}^{(l)}}{\sqrt{\mathbf{d G}^{(l) \top} \mathbf{d} \mathbf{S}^{(l)}}}, \mathbf{u}_{2}^{(l)}=\frac{\mathbb{B}^{(l)} \mathbf{d} \mathbf{S}^{(l)}}{\sqrt{\mathbf{d} \mathbf{S}^{(l) \top} \mathbb{B}_{0}^{(l)} \mathbf{d} \mathbf{S}^{(l)}}}
$$

If we simply choose $\mathbb{B}^{(l)}$ to be a constant diagonal matrix, we recover our first-order proximal gradient algorithm.
Note that an arbitrary step of updates can be implemented for these two schemes, which generally do not cost much more space due to the compact representation formula available(Nocedal and Wright, 2006). However, LSR1 generally with more steps generally performs less effectively than LBFGS due to the indefiniteness issue.

- Steps to update the stepsize. Many powerful schemes are suggested for updating the stepsize by the literature. We adopt a line search procedure with backtracking(Boyd et al., 2004). See Section 3 of Lee et al. (2014) for more details.
- Solver for the subproblem in Step 4. Step 4 breaks down to a "quadratic + nuclear norm penalty " type problem, which can be efficiently solved by first-order solvers. In our implementation, we adopt our Algorithm 1 which, in the case of quadratic smooth part, performs in the same spirit as the NNLS method by Toh and Yun (2010).

Below we run some synthetic simulation based on the multivariate regression model $(m=$ $80, N=6320$, and the ground truth is generated in the same way as the simulation in our main context) to compare the computational efficiency and statistical accuracy of first-order and second-order algorithms. Table S. 4 presents the comparison of first-order and secondorder implementations(with 1 -step limited memory SR update and 1 -step limited memory BFGS update, termed as LSR1 and LBFGS respectively). In Table S. 5 we compare BFGS approximation with different steps of storage. The "Error" column in the tables is measured by $\left\|\widehat{\mathbf{A}}-\mathbf{A}_{0}\right\|_{F}^{2}$.

| Method | Criterion | Noise |  |  |  |
| ---: | :--- | :---: | :---: | :---: | :---: |
|  |  | Gaussian | Cauchy | Lognormal |  |
| First order | Error | $9.48 \mathrm{e}-2(3.97 \mathrm{e}-3)$ | $3.21 \mathrm{e}-4(1.23 \mathrm{e}-5)$ | $1.28 \mathrm{e}-6(1.59 \mathrm{e}-7)$ |  |
|  | Rank | $5.00(0)$ | $5.00(0)$ | $5.00(0)$ |  |
|  | Iterations | 32.0 | 35.3 | 56.6 |  |
|  | Time(s) | 2.22 | 2.51 | 3.88 |  |
| LSR1 | Error | $9.48 \mathrm{e}-2(3.97 \mathrm{e}-3)$ | $3.21 \mathrm{e}-4(1.29 \mathrm{e}-5)$ | $1.48 \mathrm{e}-6(3.30 \mathrm{e}-7)$ |  |
|  | Rank | $5.00(0)$ | $5.00(0)$ | $5.00(0)$ |  |
|  | Iterations | 28.1 | 31.8 | 51.7 |  |
|  | Time(s) | 4.81 | 4.77 | 7.48 |  |
| LBFGS | Error | $9.48 \mathrm{e}-2(3.98 \mathrm{e}-3)$ | $3.20 \mathrm{e}-4(1.36 \mathrm{e}-5)$ | $1.49 \mathrm{e}-6(3.38 \mathrm{e}-7)$ |  |
|  | Rank | $5.00(0)$ | $5.00(0)$ | $5.00(0)$ |  |
|  | Iterations | 24.5 | 29.2 | 3.3 .3 |  |
|  | Time(s) | 4.19 | 4.33 |  |  |

Table S.4: Comparison of first order and second order implementation of RML

| Method | Criterion | Noise |  |  |
| :---: | :---: | :---: | :---: | :---: |
|  |  | Gaussian | Cauchy | Lognormal |
| 1-step LBFGS | Error | $9.52 \mathrm{e}-2(4.11 \mathrm{e}-3)$ | $3.19 \mathrm{e}-4(1.35 \mathrm{e}-5)$ | $1.45 \mathrm{e}-6(3.04 \mathrm{e}-7)$ |
|  | Rank | $5.00(0)$ | $5.00(0)$ | $5.00(0)$ |
|  | Iterations | 24.5 | 29.4 | 35.9 |
|  | Time(s) | 4.15 | 4.34 | 5.06 |
| 5-step LBFGS | Error | $9.52 \mathrm{e}-2(4.12 \mathrm{e}-3)$ | $3.19 \mathrm{e}-4(1.26 \mathrm{e}-5)$ | $1.41 \mathrm{e}-6(2.34 \mathrm{e}-7)$ |
|  | Rank | $5.00(0)$ | $5.00(0)$ | $5.00(0)$ |
|  | Iterations | 25.3 | 30.8 | 41.1 |
|  | Time(s) | 4.26 | 4.59 | 6.22 |
| 10-step LBFGS | Error | $9.52 \mathrm{e}-2(4.12 \mathrm{e}-3)$ | $3.19 \mathrm{e}-4(1.25 \mathrm{e}-5)$ | $1.34 \mathrm{e}-6(1.84 \mathrm{e}-7)$ |
|  | Rank | $5.00(0)$ | $5.00(0)$ | $5.00(0)$ |
|  | Iterations | 25.2 | 31.1 | 6.51 |
|  | Time(s) | 4.25 | 4.65 |  |

Table S.5: Comparison of BFGS with different steps of approximation for multivariate regression

Here is a brief summary for the tables:

- Accuracy of solutions. As we can see from the results, under the specified tolerance, all methods give quite accurate solutions. There is no clear superiority among the tested solvers. For second-order methods, both rank-1 update(LSR1) and rank-2 update(LBFGS) are able to recover the true estimand in a decent manner.
- Number of iterations. In terms of total steps of iterations, generally the rank-1 and rank-2 update(LBFGS) can reduce the steps that the algorithm requires to converge under a variety of noisy perturbations. The rank-2 update(LBFGS) appears more powerful than rank-1 update(LSR1) in accelerating the convergence.
- Time per iteration and total running time. Although the second-order schemes witness a reduction in the steps needed for convergence, they pay the cost of additional time consumption per iteration step. This is mostly due to the fact that one needs to seek the optimal solution for a sub-problem(see Algorithm 2 step 4), which again relies on first-order solvers, thus leading to more rounds of singular value decomposition. It is notable that a reduction in the number of iteration steps does lend a large improvement in sparse vector-variate problems, since solving $L_{1}$ penalization breaks down to subproblems that involves much less computation burden than matrix settings(Lee et al., 2014; Becker and Fadili, 2012). Therefore, the blessing of iteration steps and the curse of per step convergence time become a negligible trade-off in nuclear norm penalization problems.
- Sensitivity to the choice of initial curvature. First order methods can be regarded as second order methods with an extremely simple Hessian approximation
scheme(constant diagonal). However, it is in general trickier to specify an initial curvature parameter. A small specification might render the iteration divergent, while an overly large number would shrink the length of steps the algorithm can take. However, second-order programs can cleverly circumvent this issue since they are endowed with the capability of adapting themselves gradually to the appropriate scale of the Hessian structure. Therefore, while more time is required in our matrix-variate optimization settings, we do further free ourselves from the dilemma of curvature tuning and obtain more numerical stability of the solutions.
- Space of memory for the second-order implementation. The second-order limited memory implementation like BFGS can proceed with different space of memory to store the Hessian information within each step. While the statistical accuracy appears less sensitive to the Hessian memory space, more complex Hessian structure could add more numerical stability in certain cases(see the Cauchy and Lognormal column of Table S.5).


## H3: Influence of the Choice of $\alpha_{0}$ on the Estimation Error

In our theory, the number $1-\alpha_{0}$ can be regarded as the "confidence level" in the sense that our non-asymptotic bounds on the estimation error will be controlled at the optimal rate with probability close to $1-\alpha_{0}$. To see how the size of $\alpha_{0}$ influence the estimation error, we conduct some simulations in Figure S. 7 under the same setting as Example 4.1 with $m_{1}=m_{2}=40$. We show the estimation error of our estimator. In the three pictures on the left, we consider three errors, fixed sample size $n=1600$, and varying $\alpha_{0}$ belonging to $\{0.001,0.002,0.005,0.01,0.02,0.05,0.1,0.2,0.3,0.4,0.5\}$. For different random errors, when $\alpha_{0}$ is not very small, the estimation error is not so sensitive to the choice of $\alpha_{0}$. But when $\alpha_{0}$ approaches zero, the estimation error will increase rapidly, because the selected $\lambda$ will be relatively large. In the three pictures on the right, we consider three errors, three $\alpha_{0}$ 's and varying sample sizes. The $\alpha_{0}$ affects estimation error much less than sample size. In order for the optimal estimation rate to hold with high probability and good numerical performance, we suggest taking $1-\alpha_{0} \in[0.8,0.9]$.

## H4: Vertical Outliers and Leverage Points

In robust statistics, vertical outliers and leverage points are also worthy of attention. In this section, we compare the performance of different methods in these two cases.

For the vertical outliers, we consider the same setting as Example 4.1 with $m_{1}=m_{2}=$ $40, n=3200$. The random errors are sampled from $\mathcal{N}(0,0.25)+a \cdot \mathcal{N}(10,1)$, where $a$ follows the binomial distribution $\operatorname{bino}(1, \pi)$. The $\pi$ represents the average proportion of vertical outliers. We show four methods' performance at different vertical outliers proportions in the left of Figure S.8. Our method is still robust to vertical outliers.

Our method is not applicable to situations where leverage points exist. In our method, the selection of regularization parameter $\lambda$ is related to the covariates. The existence of leverage points may cause $\lambda$ to be very large, affecting the estimation error. We consider the settings in Example 4.1: $m_{1}=m_{2}=40, n=3200$, and the error is Gaussian $\mathcal{N}(0,0.25)$.


Figure S.7: The influence of the choice of $\alpha_{0}$ on estimation error of RML


Figure S.8: The influence of vertical outliers and leverage points on estimation error

Consider that $\mathbf{X}$ is generated in the following way,

$$
\mathbf{X}=\mathbf{Z}_{1}+a \cdot \mathbf{Z}_{2}
$$

where $\mathbf{Z}_{1}$ is composed of $\mathcal{N}(0,1)$ entries, $\mathbf{Z}_{2}$ is composed of $\mathcal{N}(10,1)$ entries, and $a$ follows the binomial distribution $\operatorname{bino}(1, \pi)$. The $\pi$ represents the average proportion of leverage points. We show four methods' performance at different leverage point proportions in the right of Figure S.8. As the proportion of leverage points increases, the performance of all methods collapses rapidly.

## H5: Additional Simulations for Examples 4.1

In this section, we present the additional simulation results for Examples 4.1 in the main text. We compare the performance of the proposed method with its competitors when the features come from highly correlated designs with different correlation strengths. In Example 4.1, we consider two correlation structures under sample size $n=3200$ : (i) The autoregressive covariance structure: $\operatorname{Cov}(\operatorname{vec}(\mathbf{X}))=\left(a_{i j}\right), a_{i j}=a^{|i-j|}$. (ii) The blockwise diagonal Toeplitz covariance matrix, of which each block along the diagonal is set to be

$$
\left[\begin{array}{cccccc}
1 & \frac{(p-2) a}{p-1} & \frac{(p-3) a}{p-1} & \cdots & \frac{a}{p-1} & 0 \\
\frac{(p-2) a}{p-1} & 1 & \frac{(p-2) a}{p-1} & \cdots & \frac{2 a}{p-1} & \frac{a}{p-1} \\
\vdots & & \ddots & & & \vdots \\
0 & \frac{a}{p-1} & \frac{2 a}{p-1} & \cdots & \frac{(p-2) a}{p-1} & 1
\end{array}\right]
$$

where $p=m_{1} m_{2} / 200$. Throughout, we refer $a \in(0,1)$ as the correlation strength. We show the performance of different methods under varying correlation strengths in Figure S.9. As the correlation increases, the estimation errors of all methods rise, and the proposed rank matrix lasso method still outperforms all competitors.


Figure S.9: Log Frobenius Errors for matrix regression under different correlations

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