# Matrix optimization based Euclidean embedding with outliers 

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#### Abstract

Euclidean embedding from noisy observations containing outlier errors is an important and challenging problem in statistics and machine learning. Many existing methods would struggle with outliers due to a lack of detection ability. In this paper, we propose a matrix optimization based embedding model that can produce reliable embeddings and identify the outliers jointly. We show that the estimators obtained by the proposed method satisfy a non-asymptotic risk bound, implying that the model provides a high accuracy estimator with high probability when the order of the sample size is roughly the degree of freedom up to a logarithmic factor. Moreover, we show that under some mild conditions, the proposed model also can identify the outliers without any prior information with high probability. Finally, numerical experiments demonstrate that the matrix optimization-based model can produce configurations of high quality and successfully identify outliers even for large networks.


Keywords Euclidean embedding • outliers • matrix optimizationg • low-rank matrix $\cdot$ error bound
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## 1 Introduction

Finding a complete set in a low-dimensional Euclidean space from partial noisy Euclidean distance observations, so-called embedding, is an important distance geometric problem in data science applications. In particular, when the distances are assumed to be measured in higher dimensional spaces, this leads to a typical nonlinear dimensional reduction which is widely used in statistics and machine learning. One of the biggest challenges in

[^0]embedding is that some noisy distance observations are usually contaminated with positive outlier errors. Due to the nature of applications, the outlier errors are usually much larger than the commonly assumed zero-mean measurement noises. Moreover, in many applications, the outlier errors are even larger than the true distance (e.g., the Non-Line-Of-Sight (NLOS) errors from wireless sensor network localization [32,37,16]), which is the main reason for the distance based methods often fail in practical embedding applications.

Needless to say, it is important to mitigate outlier errors from the observation distances for embedding. In literature, for some applications (e.g., wireless communication), one way to mitigate outlier propagation is to develop some methods to identify outlier errors through prior information, such as the outlier distribution and physical characteristics of networks [41].35]. For the overview of various outlier identification techniques and optimization methods in wireless communication applications, see the nice survey [16]. However, in most applications, the prior information of outliers is either technically non-available or costly to obtain due to hardware limitations. Therefore, it is even more crucial to identify and mitigate the outlier propagation from the observed distance data without prior information. To this end, different matrix optimization models are proposed and become popular in applications. For instance, in wireless sensor localization, different semidefinite programming (SDP) based methods are proposed by [7, 40, 48], and numerical experiments demonstrate that SDP based models can provide descent estimations even without prior information on outliers for some small-scale applications. More recently, based on the concept of the Euclidean distance matrix (EDM), a new matrix optimization model for the outlier mitigation has been proposed in [10] (see Section 2 for details). Numerical tests on both simulations and real-world applications show that the EDM based method proposed in [10] can produce high quality embeddings without prior information even for large-scale networks. Numerical experiments show that one of the main advantages of the EDM model [10] comparing with the existing SDP approaches is that the EDM model usually is able to identify outlier errors index sets (i.e., the index sets of observation distances which contain outlier errors). However, there is no theoretical guarantee on the outlier detection ability provided in [10]. The main purpose of this paper is to study the statistical performance analysis of the EDM based embedding with outliers by establishing the recovery error bounds and the embedding dimension and outlier detection guarantee.

In general, the EDM based embedding model proposed in [10] belongs the category of low-rank matrix approximation problems [27]14], which had many exciting developments recently and attracted much attention from optimization and machine learning communities. More precisely, in principle, the proposed EDM model is in line with the general framework of robust principal component analysis (Robust PCA) [6]3], i.e., estimating an unknown low-rank matrix $\bar{X} \in \mathbb{R}^{m \times n}$ from a collection of partially noisy observed elements $\left.\widetilde{X}_{i j}=\bar{X}_{i j}+\bar{S}_{i j}+\right\rceil_{i j}$, $(i, j) \in \Omega \subseteq\{1, \ldots, m\} \times\{1, \ldots, n\}$, where $\bar{S}$ is a sparse matrix consisting of outliers, 7 represents the random noise and $\Omega$ is the observation index subset.

Enlightened by the previous tremendous success of the convex matrix optimization approaches in low-rank matrix completion [5]4], Chandrasekaran et al. [6] first study the Robust PCA for the case that the elements $\widetilde{X}$ are fully observed and without noise. In particular, Chandrasekaran et al. [6] showed that under the rank-sparsity incoherent property, the unknown true low-rank matrix $\bar{X}$ and sparse outlier matrix $\bar{S}$ can be recovered exactly based on the convex "nuclear norm plus $l_{1}$-norm" approach. For the general setting with missing observations, by employing the previous developed probability analysis techniques for exact matrix completion problems [5] 34 15], Candès et al. [3] provided the probabilistic guarantees for exact recovery of the convex "nuclear norm plus $l_{1}$-norm" approach for Robust PCA. Later a sharper probabilistic exact recovery guarantee was established by Chen et al. [9] focused on high-dimensional statistical settings (i.e., the sample size is smaller than mn ).

For the more realistic noisy setting, Zhou et al. [49] proposed a convex "nuclear norm plus $l_{1}$-norm" constrained matrix optimization model and studied its statistical performance guarantees, and later the nuclear norm plus $l_{1}$-norm penalized formulation was studied by Hsu et al. [19]. Based on the unified restricted strong convexity (RSC) framework introduced by [31,30], Agarwal et al. [1] obtain a sharper statistical error bound for the nuclear norm plus $l_{1}$-norm penalized model. However, in prior studies [49, 19, 1] on the noisy robust PCA, the performance guarantee results are all based on the full observation assumption, which may not be practical in applications. In [45], Wong and Lee established an estimation error bound for the noisy robust PCA, with an
assumption that the number of observed entries is in the order of $m n$. However, this also may not be useful for high-dimensional applications, since the sample sizes there are usually much smaller. For high-dimensional settings, under the boundedness assumptions on the true low-rank matrix $\bar{X}$ and sparse outlier matrix $\bar{S}$, Klopp, et al. [22] derived a statistical estimation error bounds for the $l_{\infty}$ constrained convex program based estimators, in which the minimal sample size required for a faithful estimation is roughly on the order of $\max \{m, n\} r \log (n)$. Recently, Chen, et al. [8], improved and derived a near-optimal statistical guarantee of the convex nuclear norm plus $l_{1}$-norm penalized model for the Robust PCA by building up the connection between the convex estimations and an auxiliary nonconvex optimization algorithm.

The results mentioned above are all about the "classical" robust PCA, in the sense that for the proposed models there is no "hard-constraints", e.g., the noisy correlation matrix recovery (i.e., a positive semidefinite matrix whose diagonal elements are all ones) and the EDM estimation considered in this paper. However, these "hardconstraints" are usually crucial and must be satisfied in the convex estimation models in many applications e.g., the EDM embedding. Consequently, the results obtained in [45, 22|,8] have become inadequate in these applications. For the correlation matrix estimation problem, Wu [46] first studied the probabilistic guarantees of the Robust PCA with "hard-constraints" for both noiseless and noisy cases. The main techniques employed in [46] are the unified restricted strong convexity (RSC) framework introduced by [31,30] and a matrix Bernstein inequality (cf. e.g., [42]), which are wildly used in the study of statistical performance guarantees of convex models in matrix completion problem (e.g., [15|,34,22, 30, 21, 29, 28|) and the EDM embedding problem without outliers [11].

In order to establish the theoretical performance analysis of the convex matrix optimization model for EDM embedding with outliers, we first adopt the error bound analysis approach introduced in [46] to EDM embedding with outliers and obtain the statistical guarantee of the convex estimation model. Furthermore, based on the resulting error bound results, we show that under some wild conditions, with high probability, the convex EDM estimator will recover the true unknown embedding dimension. Simultaneously, we also show that the outlier estimator obtained by the convex matrix optimization model will recover the index set of the support set of the unknown outliers with the same probability. Finally, we verify the proposed theoretical results by numerical experiments.

The remaining parts of this paper are organized as follows. We briefly introduce the matrix optimization based EDM embedding with outliers model originally proposed in [10]. Section 3 contains the statistical recovery error bounds for the EDM embedding EDM model. In Section 4, we establish the probability recovery guarantee of the embedding dimensionality and outlier detection. We verify the theoretical results obtained in Sections 3 and 4 through numerical examples in Section 5 We conclude the paper in Section 6

Below are some common notations to be used in this paper:

- For any $Z \in \mathbb{R}^{m \times n}$, we denote by $Z_{i j}$ the $(i, j)$-th entry of $Z$.
- We use " $\circ$ " to denote the Hadamard product between matrices, i.e., for any two matrices $X$ and $Y$ in $\mathbb{R}^{m \times n}$ the $(i, j)$-th entry of $Z:=X \circ Y \in \mathbb{R}^{m \times n}$ is $Z_{i j}=X_{i j} Y_{i j}$.
- For any $Z \in \mathbb{R}^{m \times n}$, we use $Z^{1 / 2} \in \mathbb{R}^{m \times n}$ to denote the $m \times n$ matrix whose $(i, j)$-th entry is $Z_{i j}^{1 / 2}$.
- Let $\mathbf{1} \in \mathbb{R}^{n}$ be the vector whose elements are all ones. Denote the $n \times n$ identity matrix by $I$ and the centering matrix by

$$
\begin{equation*}
J:=I-\mathbf{1 1}^{T} / n . \tag{1}
\end{equation*}
$$

- For a given $Z \in \mathbb{S}^{n}$, we use $\lambda_{1}(Z) \geq \ldots \geq \lambda_{n}(Z)$ to denote the eigenvalues of $Z$ (all real and counting multiplicity) arranging in non-increasing order and use $\lambda(Z)$ to denote the vector of the ordered eigenvalues of $Z$. Let $\mathbb{O}^{n}$ be the set of all $n \times n$ orthogonal matrices.
- Let $\operatorname{sgn}(\cdot): \mathbb{R}^{m \times n} \rightarrow \mathbb{R}^{m \times n}$ be the (component-wise) sign function, i.e., for any $Z \in \mathbb{R}^{m \times n}, \operatorname{sgn}(Z)_{i j}=1$ if $Z_{i j}>0$ and $\operatorname{sgn}(Z)_{i j}=0$ otherwise.
- For any $z \in \mathbb{R}^{n}$, we use $\operatorname{Diag}(z)$ to denote an $n \times n$ diagonal matrix with $z$ on the main diagonal. Meanwhile, for any $Z \in \mathbb{R}^{n \times n}$, $\operatorname{diag}(Z)$ denotes the column vector consisting of all the diagonal entries of $Z$ being arranged from the first to the last.


## 2 The EDM based embedding with outliers

Let $p_{i}, i=1, \ldots, n$ be $n$ points in a $r$-dimensional subspace. For each $i, j \in\{1, \ldots, n\}$, denote $\bar{d}_{i j} \geq 0$ the distance between the $i$-th and $j$-th points on the $r$-dimensional subspace. A $n \times n$ matrix $\bar{D}$ is called Euclidean distance matrix (EDM) if $\bar{D}_{i j}=\left(\bar{d}_{i j}\right)^{2}$ for $i, j=1, \ldots, n$. An alternative definition of EDM that does not involve any embedding points $\left\{p_{i}\right\}$ can be described as follows. Let $\mathbb{S}_{h}^{n}$ be the hollow subspace of $\mathbb{S}^{n}$, i.e., $\mathbb{S}_{h}^{n}:=\left\{X \in \mathbb{S}^{n} \mid \operatorname{diag}(X)=0\right\}$. Define the almost positive semidefinite cone $\mathbb{K}_{+}^{n}$ by

$$
\begin{equation*}
\mathbb{K}_{+}^{n}:=\left\{A \in \mathbb{S}^{n} \mid x^{T} A x \geq 0 \forall x \in \mathbf{1}^{\perp}\right\}=\left\{A \in \mathbb{S}^{n} \mid J A J \succeq 0\right\} \tag{2}
\end{equation*}
$$

where $\mathbf{1}^{\perp}:=\left\{x \in \mathbb{R}^{n} \mid \mathbf{1}^{T} x=0\right\}$. It is well-known [36 47] that $D \in \mathbb{S}^{n}$ is EDM if and only if $-D \in \mathbb{S}_{h}^{n} \cap \mathbb{K}_{+}^{n}$. Moreover, the embedding dimension is determined by the rank of the doubly centered matrix $J \bar{D} J$, i.e., $r=$ $\operatorname{rank}(J \bar{D} J)$. Given a true EDM $\bar{D}$, since $-J \bar{D} J$ is positive semidefinite, its spectral decomposition can be written as

$$
-\frac{1}{2} J \bar{D} J=P \operatorname{Diag}\left(\lambda_{1}, \ldots, \lambda_{n}\right) P^{T},
$$

where $P \in \mathbb{O}^{n}$ and $\lambda_{1} \geq \lambda_{2} \geq \cdots \geq \lambda_{n} \geq 0$ are the eigenvalues in nonincreasing order. Let $P_{1}$ be the submatrix consisting of the first $r$ columns (eigenvectors) in $P$. One set of the embedding points are

$$
\left(\begin{array}{l}
p_{1}^{T}  \tag{3}\\
\vdots \\
p_{n}^{T}
\end{array}\right)=P_{1} \operatorname{Diag}\left(\sqrt{\lambda_{1}}, \ldots, \sqrt{\lambda_{r}}\right) .
$$

In order to find a set of relative embedding points $\left\{p_{i}\right\}$, we are interesting in estimating the ture EDM $\bar{D}$ from the partial noisy observation distances $\tilde{d}_{i j}$. The basic noisy model takes the following form

$$
\begin{equation*}
\tilde{d}_{i j}=\bar{d}_{i j}+\eta \xi_{i j}, \quad i, j \in\{1, \ldots, n\} . \tag{4}
\end{equation*}
$$

where $\xi_{i j}$ are i.i.d. noise errors with $\mathbb{E}\left(\xi_{i j}\right)=0$ and $\mathbb{E}\left(\xi_{i j}^{2}\right)=1, \eta>0$ is a noise magnitude control factor.
Unlike the standard zero-mean noise assumption, we are attractive to the case where the distance measurements $\tilde{d}_{i j}$ are also contaminated with the errors arising from outliers, which usually have significant positive biases and cause the measured distances $\tilde{d}_{i j}$ significantly diverging from actual values $\bar{d}_{i j}$. The errors from outliers frequently appear in many applications such as the Non-Line-Of-Sight (NLOS) errors from wireless sensor network localization [32,37, 16], the errors arising from outliers in manifold learning and others [18, 23].

We use $\bar{s} \in \mathbb{S}^{n}$ to represent the outlier errors, whose elements are either zero or positive accordingly. Therefore, the basic noisy model (4) then takes the following form

$$
\begin{equation*}
0 \leq \tilde{d}_{i j}=\bar{d}_{i j}+\bar{s}_{i j}+\eta \xi_{i j}, \quad i, j \in\{1, \ldots, n\} . \tag{5}
\end{equation*}
$$

For notational simplicity, we define the outlier matrix $\bar{S} \in \mathbb{S}^{n}$ as follows

$$
\begin{equation*}
\bar{S}_{i j}:=\bar{s}_{i j}\left(\bar{s}_{i j}+2 \bar{d}_{i j}\right) \geq 0, \quad i, j \in\{1, \ldots, n\} . \tag{6}
\end{equation*}
$$

Then, by (5], we know that

$$
\begin{equation*}
\widetilde{D}_{i j}=\tilde{d}_{i j}^{2}=\left(\bar{d}_{i j}+\bar{s}_{i j}+\eta \xi_{i j}\right)^{2}=\bar{D}_{i j}+\bar{S}_{i j}+ד_{i j}, \quad i, j \in\{1, \ldots, n\}, \tag{7}
\end{equation*}
$$

where $\widetilde{D}$ and $\bar{D} \in \mathbb{S}^{n}$ are the observation and unknown true EDMs, whose $(i, j)$-element are $\tilde{d}_{i j}^{2}$ and $\vec{d}_{i j}^{2}$, respectively, and $\rceil \in \mathbb{S}^{n}$ is defined by

$$
\begin{equation*}
T_{i j}:=2\left(\bar{d}_{i j}+\bar{s}_{i j}\right) \eta \xi_{i j}+\eta^{2} \xi_{i j}^{2}, \quad i, j \in\{1, \ldots, n\} \tag{8}
\end{equation*}
$$

Without loss of generality, we always assume that the magnitude of measurement error $\left|\eta \xi_{i j}\right|$ is strictly smaller than the true distance $d_{i j}$ for each $i, j \in\{1, \ldots, n\}$. Thus, since $\tilde{d}_{i j} \geq 0$, we know that for any $i, j \in\{1, \ldots, n\}$,

$$
\bar{s}_{i j}=0 \quad \Longleftrightarrow \quad \bar{S}_{i j}=0
$$

For each $i \in\{1, \ldots, n\}$, we use $\mathbf{e}_{i} \in \mathbb{R}^{n}$ to denote the $i$-th canonical basis of $\mathbb{R}^{n}$. Let $\mathbb{S}_{h}^{n} \subseteq \mathbb{S}^{n}$ be the hollow space, i.e.,

$$
\mathbb{S}_{h}^{n}:=\left\{Z \in \mathbb{S}^{n} \mid Z_{i i}=0, i=1, \ldots, n\right\}
$$

whose dimension equals to $d_{\mathbb{S}_{h}^{n}} \equiv n(n-1) / 2$. Let $\left\{\frac{1}{2}\left(\mathbf{e}_{i} \mathbf{e}_{j}^{T}+\mathbf{e}_{j} \mathbf{e}_{i}^{T}\right)\right\}_{1 \leq i<j \leq n}$ be the standard basis matrices of $\mathbb{S}_{h}^{n}$.
For the given observation set $\Omega$, let $\left\{X_{1}, \ldots, X_{m}\right\}$ with $m:=|\Omega|$ be the numbered sampled basis matrices from the standard basis matrices set $\left\{\frac{1}{2}\left(\mathbf{e}_{i} \mathbf{e}_{j}^{T}+\mathbf{e}_{j} \mathbf{e}_{i}^{T}\right)\right\}_{1 \leq i<j \leq n}$. Therefore, the corresponding observation operator $\mathscr{O}_{\Omega}: \mathbb{S}^{n} \rightarrow \mathbb{R}^{m}$ can be written as

$$
\begin{equation*}
\mathscr{O}_{\Omega}(A):=\left(\left\langle X_{1}, A\right\rangle, \ldots,\left\langle X_{m}, A\right\rangle\right)^{T} \in \mathbb{R}^{m}, \quad A \in \mathbb{S}^{n} . \tag{9}
\end{equation*}
$$

That is, $\mathscr{O}_{\Omega}(A)$ samples all the elements $A_{i j}$ specified by $(i, j) \in \Omega$. Let $\mathscr{O}_{\Omega}^{*}: \mathbb{R}^{m} \rightarrow \mathbb{S}^{n}$ be its adjoint, i.e., $\mathscr{O}_{\Omega}^{*}(\mathbf{z})=$ $\sum_{l=1}^{m} \mathbf{z}_{l} X_{l}, \mathbf{z} \in \mathbb{R}^{m}$. Then, we further define the observation vector

$$
\begin{equation*}
\mathbf{y}:=\mathscr{O}_{\Omega}(\widetilde{D}) \in \mathbb{R}^{m} \tag{10}
\end{equation*}
$$

Finally, under the assumption that $\bar{S}$ is sparse (i.e., the cardinality of nonzero elements of $\bar{S}$ is small), we may estimate the unknown matrices $\bar{D}$ and $\bar{S}$ by solving the following nonconvex optimization model:

$$
\begin{array}{ll}
\min & \frac{1}{2 m}\left\|\mathbf{y}-\mathscr{O}_{\Omega}(D+S)\right\|^{2}+\rho\|S\|_{0} \\
\text { s.t. } & D \in \mathbb{S}_{h}^{n}, \quad-D \in \mathbb{K}_{+}^{n}, \quad \operatorname{rank}(-J D J) \leq r,  \tag{11}\\
& S \in \mathbb{S}^{n}, \quad S \geq 0
\end{array}
$$

Since the rank constraint and the zero norm $\|\cdot\|_{0}$ are computational intractable, we may consider the following convex relaxated matrix optimization problem

$$
\begin{array}{ll}
\min & \Phi_{\rho_{D}, \rho_{S}}(D, S):=\frac{1}{2 m}\left\|\mathbf{y}-\mathscr{O}_{\Omega}(D+S)\right\|^{2}+\rho_{D}\langle I-\widetilde{F},-J D J\rangle+\rho_{S}\langle E-\widetilde{G}, S\rangle  \tag{12}\\
\text { s.t. } & D \in \mathbb{S}_{h}^{n}, \quad-D \in \mathbb{K}_{+}^{n}, \quad S \geq 0
\end{array}
$$

where $\rho_{D}$ and $\rho_{S}$ are two given positive parameters, $E \in \mathbb{S}^{n}$ is the matrix whose elements are all ones, and $\widetilde{F}$ and $\widetilde{G}$ are two given symmetric matrices. In particular, when both $\widetilde{F}$ and $\widetilde{G}$ vanish, the model (12) reduces to the following convex nuclear norm $l_{1}$-minimization EDM matrix optimization problem:

$$
\begin{array}{ll}
\min & \frac{1}{2 m}\left\|\mathbf{y}-\mathscr{O}_{\Omega}(D+S)\right\|^{2}+\rho_{D}\langle I,-J D J\rangle+\rho_{S}\langle E, S\rangle  \tag{13}\\
\text { s.t. } & D \in \mathbb{S}_{h}^{n}, \quad-D \in \mathbb{K}_{+}^{n}, \quad S \geq 0
\end{array}
$$

We use ( $\widehat{D}_{m}, \widehat{S}_{m}$ ) to denote an optimal solution of the above convex model (12), and in later discussions, we often drop the subscript " $m$ ", when the dependence of $(\widehat{D}, \widehat{S}$ ) on the sample size $m$ is clear from the context. Furthermore, we use $(\bar{D}, \bar{S})$ to denote the unknown true EDM and outlier matrix.

In this paper, we choose the symmetric matrices $\widetilde{F} \in \mathbb{S}^{n}$ and $\widetilde{G} \in \mathbb{S}^{n}$ in the objective function of (12) by following the suggestions [29, (25) and (26)] (see also [46, Chapter 5.3]). Suppose that the EDM $\widetilde{D}$ and the non-negative matrix $\widetilde{S}$ are the given initial estimators (e.g., the estimators obtained by 133 ). Define $\widetilde{F} \in \mathbb{S}^{n}$ by

$$
\begin{equation*}
\widetilde{F}=F(-J \widetilde{D} J), \tag{14}
\end{equation*}
$$

where $F: \mathbb{S}^{n} \rightarrow \mathbb{S}^{n}$ is the spectral operator [12, 13] associated with the symmetric function (cf. [12, Definition 1]) $f: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ defined by

$$
f_{i}(x)=\left\{\begin{array}{ll}
\phi\left(x_{i} / \max _{l}\left\{x_{l}\right\}\right) & \text { if } x \neq 0,  \tag{15}\\
0 & \text { otherwise },
\end{array} \quad x \in \mathbb{R}^{n}\right.
$$

with the scalar function $\phi: \mathbb{R} \rightarrow \mathbb{R}$ takes the form

$$
\begin{equation*}
\phi(t)=\left(1+\varepsilon^{\tau}\right) \frac{t^{\tau}}{t^{\tau}+\varepsilon^{\tau}}, \quad t \in \mathbb{R} \tag{16}
\end{equation*}
$$

for some $\varepsilon>0$ and $\tau>0$. Meanwhile, we define the symmetric matrix $\widetilde{G} \in \mathbb{S}^{n}$ with respect to $\widetilde{S}$ by

$$
\widetilde{G}_{i j}=\left\{\begin{array}{ll}
\phi\left(\widetilde{S}_{i j} / \max _{k, l}\left\{\widetilde{S}_{k l}\right\}\right) & \text { if } \widetilde{S} \neq 0,  \tag{17}\\
0 & \text { otherwise },
\end{array} \quad(i, j) \in\{1, \ldots, n\}\right.
$$

where $\phi: \mathbb{R} \rightarrow \mathbb{R}$ is the scalar function defined by 16 .
Throughout this paper, the following condition is assumed to hold, which ensures the existent of the optimal solution of (12).

Assumption 1 There exists constant $\bar{\rho}_{D}>0$ and $\bar{\rho}_{S}>0$ such that for any $\alpha$, the level set

$$
L_{\bar{\rho}_{D}, \bar{\rho}_{S}}(\alpha):=\left\{(D, S) \in \mathbb{S}^{n} \mid \Phi_{\bar{\rho}_{D}, \bar{\rho}_{S}}(D, S) \leq \alpha, D \in \mathbb{S}_{h}^{n},-D \in \mathbb{K}_{+}^{n}, S \geq 0\right\}
$$

of (12) is closed and bounded.
It is worth to note that for any $0 \leq t \leq 1, \phi(t) \in[0,1]$. Thus, it is easy to check that for any given initial EDM estimator $\widetilde{D}$ and outlier matrix estimator $\widetilde{S}$, the symmetric matrices $\widetilde{F}$ and $\widetilde{G}$ satisfy $\langle I-\widetilde{F},-J D J\rangle \geq 0$ and $\langle E-\widetilde{G}, S\rangle \geq 0$. Therefore, we know that for any $\rho_{D} \geq \bar{\rho}_{D}>0$ and $\rho_{S} \geq \bar{\rho}_{S}>0, L_{\rho_{D}, \rho_{S}}(\alpha) \subseteq L_{\bar{\rho}_{D}, \bar{\rho}_{S}}(\alpha)$ for any $\alpha$. This yields that under Assumption 1, for any $\rho_{D} \geq \bar{\rho}_{D}>0$ and $\rho_{S} \geq \bar{\rho}_{S}>0$, the optimal solution of (12) exists. Moreover, under Assumption 1, we know that there exist two positive constants $b_{D}$ and $b_{S}$ such that the optimal solution $(\widehat{D}, \widehat{S})$ of 12 for any $\rho_{D} \geq \bar{\rho}_{D}$ and $\rho_{S} \geq \bar{\rho}_{S}$ and the unkown true EDM and outlier matrix $(\bar{D}, \bar{S})$ satisfy

$$
\|\widehat{D}\|_{\infty} \leq b_{D}, \quad\|\bar{D}\|_{\infty} \leq b_{D}, \quad\|\widehat{S}\|_{\infty} \leq b_{S} \quad \text { and } \quad\|\bar{S}\|_{\infty} \leq b_{S}
$$

## 3 Recovery error bounds

In this section, we aim to derive a recovery error bound for the proposed matrix optimization based EDM embedding model with outliers (12]. Here we adopt the approach introduced by Wu [46] for studying recovery error bounds of the noisy low-rank and sparse matrix decomposition with fixed basis. Essentially, the proofs are in line with the well-studied unified framework introduced by Negahban et al. [31] for high-dimensional analysis of M-estimators with decomposable regularizers, which is used frequently in the study of noisy matrix completion [30, 21, 28, 29 46, 11]. For the sake of completion, we include the detailed proofs in Appendix.

Recall that $\bar{D}$ is the unknown true EDM matrix. Suppose that the positive semidefinite matrix $-J \bar{D} J$ has the eigenvalue decomposition

$$
\begin{equation*}
-J \bar{D} J=\bar{P} \bar{\Lambda} \bar{P}^{T} \tag{18}
\end{equation*}
$$

where $\bar{\Lambda} \in \mathbb{S}^{n}$ is a diagonal matrix whose diagonal elements are the eigenvalues of $-J \bar{D} J$ arranged in the nonincreasing order, and $\bar{P}=\left[\bar{P}_{1} \bar{P}_{2}\right] \in \mathbb{O}^{n}$ with $\bar{P}_{1} \in \mathbb{R}^{n \times r}$ and $\bar{P}_{2} \in \mathbb{R}^{n \times(n-r)}$. We define the generalized geometric
center subspace in $\mathbb{S}^{n}$ with respect to $\bar{P}_{1}$ by $\mathbb{T}:=\left\{Y \in \mathbb{S}^{n} \mid Y \bar{P}_{1}=0\right\}$. Let $\mathbb{T}^{\perp}$ be its orthogonal subspace. Then, the orthogonal projectors to the two subspaces can hence be calculated respectively by

$$
\begin{equation*}
\mathscr{P}_{\mathbb{T}}(A):=\bar{P}_{2} \bar{P}_{2}^{T} A \bar{P}_{2} \bar{P}_{2}^{T} \quad \text { and } \quad \mathscr{P}_{\mathbb{T}^{\perp}}(A):=\bar{P}_{1} \bar{P}_{1}^{T} A+A \bar{P}_{1} \bar{P}_{1}^{T}-\bar{P}_{1} \bar{P}_{1}^{T} A \bar{P}_{1} \bar{P}_{1}^{T}, \quad A \in \mathbb{S}^{n} \tag{19}
\end{equation*}
$$

Moreover, we have the following orthogonal decomposition

$$
\begin{equation*}
A=\mathscr{P}_{\mathbb{T}}(A)+\mathscr{P}_{\mathbb{T}^{\perp}}(A) \quad \text { and } \quad\left\langle\mathscr{P}_{\mathbb{T}}(A), \mathscr{P}_{\mathbb{T}^{\perp}}(B)\right\rangle=0 \quad \forall A, B \in \mathbb{S}^{n} \tag{20}
\end{equation*}
$$

It then follows from the definition of $\mathscr{P}_{\mathbb{T}}$ that for any $A \in \mathbb{S}^{n}, \mathscr{P}_{\mathbb{T}^{\perp}}(A)=\bar{P}_{1} \bar{P}_{1}^{T} A+\bar{P}_{2} \bar{P}_{2}^{T} A \bar{P}_{1} \bar{P}_{1}^{T}$, which implies that $\operatorname{rank}\left(\mathscr{P}_{\mathbb{T}^{\perp}}(A)\right) \leq 2 r$. This yields that for any $A \in \mathbb{S}^{n}$,

$$
\begin{equation*}
\left\|\mathscr{P}_{\mathbb{T}^{\perp}}(A)\right\|_{*} \leq \sqrt{2 r}\|A\| . \tag{21}
\end{equation*}
$$

For any given $S \in \mathbb{S}^{n}$, we use $\operatorname{supp}_{S}$ to denote the index set of the support of $S$, i.e.,

$$
\operatorname{supp}(S):=\left\{(i, j) \mid S_{i j} \neq 0, i, j \in\{1, \ldots, n\}\right\}
$$

Suppose that the unknown matrix $\bar{S}$ has $k$ nonzero entries, i.e., $\|\bar{S}\|_{0}=|\operatorname{supp}(\bar{S})|=k$. The tangent subspace $\mathbb{L}$ with respect to the $k$-sparse subset $\left\{S \in \mathbb{S}^{n} \mid\|S\|_{0} \leq k\right\}$ at $\bar{S}$ then takes the form

$$
\mathbb{L}:=\left\{S \in \mathbb{S}^{n} \mid \operatorname{supp}(S) \subseteq \operatorname{supp}(\bar{S})\right\} .
$$

Denote the orthogonal complement of $\mathbb{L}$ by $\mathbb{L}^{\perp}$. Let $\mathscr{P}_{\mathbb{L}}$ and $\mathscr{P}_{\mathbb{L}^{\perp}}$ be the corresponding orthogonal projections onto $\mathbb{L}$ and $\mathbb{L}^{\perp}$. Then, we have the following decomposition

$$
\begin{equation*}
B=\mathscr{P}_{\mathbb{L}}(B)+\mathscr{P}_{\mathbb{L}^{\perp}}(B) \quad \text { and } \quad\left\langle\mathscr{P}_{\mathbb{L}}(A), \mathscr{P}_{\mathbb{L}^{\perp}}(B)\right\rangle=0 \quad \forall A, B \in \mathbb{S}^{n} . \tag{22}
\end{equation*}
$$

Moreover, for any $B \in \mathbb{S}^{n}$, since $\left\|\mathscr{P}_{\mathbb{L}}(B)\right\|_{0} \leq k$, we have that

$$
\begin{equation*}
\left\|\mathscr{P}_{\mathbb{L}}(B)\right\|_{1} \leq \sqrt{k}\|B\| . \tag{23}
\end{equation*}
$$

Define

$$
\begin{equation*}
a_{D}:=\frac{1}{\sqrt{r}}\left\|\bar{P}_{1} \bar{P}_{1}^{T}-\widetilde{F}\right\| \quad \text { and } \quad a_{S}:=\frac{1}{\sqrt{k}}\|\operatorname{sgn}(\bar{S})-\widetilde{G}\| . \tag{24}
\end{equation*}
$$

It is also easy to verify the following result.
Lemma 1 For any $D \in \mathbb{S}_{h}^{n}$, we have $D-J D J=\frac{1}{2}\left(\operatorname{diag}(-J D J) \mathbf{1}^{T}+\mathbf{1} \operatorname{diag}(-J D J)^{T}\right)$.
The following result represents the first important step to derive our error bounds of the convex model (12).
Proposition 1 Let $(\widehat{D}, \widehat{S})$ and $(\bar{D}, \bar{S})$ be an optimal solution of (12) and the underground true EDM and outlier matrices, respectively. Let $\kappa_{D}>1$ and $\kappa_{S}>1$ be given. Suppose that the parameters $\rho_{D}$ and $\rho_{S}$ satisfy

$$
\begin{equation*}
\rho_{D} \geq \frac{\kappa_{D}}{m}\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{2} \quad \text { and } \quad \rho_{S} \geq \frac{\kappa_{S}}{m}\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{\infty} \tag{25}
\end{equation*}
$$

where $\left.\zeta:=\mathscr{O}_{\Omega}( \rceil\right) \in \mathbb{R}^{m}$ and $\rceil \in \mathbb{S}^{n}$ is given by (8), then we have

$$
\begin{equation*}
\frac{1}{2 m}\left\|\mathscr{O}_{\Omega}(\widehat{D}-\bar{D})+\mathscr{O}_{\Omega}(\widehat{S}-\bar{S})\right\|^{2} \leq \rho_{D} \sqrt{r}\left(a_{D}+\frac{2 \sqrt{2}}{\kappa_{D}}\right)\|\widehat{D}-\bar{D}\|+\rho_{S} \sqrt{k}\left(a_{S}+\frac{1}{\kappa_{S}}\right)\|\widehat{S}-\bar{S}\| \tag{26}
\end{equation*}
$$

and

$$
\left\{\begin{array}{l}
\|\widehat{D}-\bar{D}\|_{*} \leq \frac{\kappa_{D}}{\kappa_{D}-1}\left(\left(a_{D}+2 \sqrt{2}\right) \sqrt{r}\|\widehat{D}-\bar{D}\|+\frac{\rho_{S}}{\rho_{D}}\left(a_{S}+\frac{1}{\kappa_{S}}\right) \sqrt{k}\|\widehat{S}-\bar{S}\|\right)  \tag{27}\\
\|\widehat{S}-\bar{S}\|_{1} \leq \frac{\kappa_{S}}{\kappa_{S}-1}\left(\frac{\rho_{D}}{\rho_{S}}\left(a_{D}+\frac{2 \sqrt{2}}{\kappa_{D}}\right) \sqrt{r}\|\widehat{D}-\bar{D}\|+\left(a_{S}+1\right) \sqrt{k}\|\widehat{S}-\bar{S}\|\right)
\end{array}\right.
$$

where $a_{D}$ and $a_{S}$ are given by (24).

Proof By (10) and (7), we know that for any $D$ and $S \in \mathbb{S}^{n}$,

$$
\begin{align*}
\frac{1}{2 m}\left\|\mathbf{y}-\mathscr{O}_{\Omega}(D+S)\right\|^{2} & =\frac{1}{2 m}\left\|\mathscr{O}_{\Omega}(\bar{D}+\bar{S})+\zeta-\mathscr{O}_{\Omega}(D+S)\right\|^{2} \\
& =\frac{1}{2 m}\left\|\mathscr{O}_{\Omega}(D-\bar{D})+\mathscr{O}_{\Omega}(S-\bar{S})\right\|^{2}-\frac{1}{m}\left\langle D-\bar{D}+S-\bar{S}, \mathscr{O}_{\Omega}^{*}(\zeta)\right\rangle+\frac{1}{2 m}\|\zeta\|^{2} \tag{28}
\end{align*}
$$

Since $(\bar{D}, \bar{S})$ is a feasible solution of $(12)$, we know from the optimality of $(\widehat{D}, \widehat{S})$ that

$$
\begin{align*}
\frac{1}{2 m}\left\|\mathscr{O}_{\Omega}(\widehat{D}-\bar{D})+\mathscr{O}_{\Omega}(\widehat{S}-\bar{S})\right\|^{2} \leq & \frac{1}{m}\left\langle\mathscr{O}_{\Omega}^{*}(\zeta), \widehat{D}-\bar{D}+\widehat{S}-\bar{S}\right\rangle-\rho_{D}(\langle I,-J(\widehat{D}-\bar{D}) J\rangle-\langle\widetilde{F},-J(\widehat{D}-\bar{D}) J\rangle) \\
& -\rho_{S}\left(\|\widehat{S}\|_{1}-\|\bar{S}\|_{1}-\langle\widetilde{G}, \widehat{S}-\bar{S}\rangle\right) \tag{29}
\end{align*}
$$

By the Hölder inequality, we know that the first term of the right hand side of 29) satisfies

$$
\frac{1}{m}\left\langle\mathscr{O}_{\Omega}^{*}(\zeta), \widehat{D}-\bar{D}+\widehat{S}-\bar{S}\right\rangle \leq \frac{1}{m}\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{2}\|\widehat{D}-\bar{D}\|_{*}+\frac{1}{m}\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{\infty}\|\widehat{S}-\bar{S}\|_{1}
$$

Since $\|\widehat{D}-\bar{D}\|_{*}=\|\widehat{D}-\bar{D}-J(\widehat{D}-\bar{D}) J+J(\widehat{D}-\bar{D}) J\|_{*}$, we know from Lemma 1 that

$$
\|\widehat{D}-\bar{D}\|_{*} \leq\|\widehat{D}-\bar{D}-J(\widehat{D}-\bar{D}) J\|_{*}+\|J(\widehat{D}-\bar{D}) J\|_{*} \leq \sqrt{2}\|\widehat{D}-\bar{D}-J(\widehat{D}-\bar{D}) J\|+\|J(\widehat{D}-\bar{D}) J\|_{*}
$$

Moreover, since $\langle J(\widehat{D}-\bar{D}) J, \widehat{D}-\bar{D}-J(\widehat{D}-\bar{D}) J\rangle=0$, we have $\|\widehat{D}-\bar{D}\|^{2}=\|\widehat{D}-\bar{D}-J(\widehat{D}-\bar{D}) J\|^{2}+\| J(\widehat{D}-$ $\bar{D}) J \|^{2}$, which implies that $\|\widehat{D}-\bar{D}\|_{*} \leq \sqrt{2}\|\widehat{D}-\bar{D}\|+\|-J(\widehat{D}-\bar{D}) J\|_{*}$. Thus, since $\|J(\widehat{D}-\bar{D}) J\| \leq\|\widehat{D}-\bar{D}\|$, by (25) and 21, we know that

$$
\begin{aligned}
\frac{1}{m}\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{2}\|\widehat{D}-\bar{D}\|_{*} & \leq \frac{1}{m}\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{2}\left(\sqrt{2}\|\widehat{D}-\bar{D}\|+\|-J(\widehat{D}-\bar{D}) J\|_{*}\right) \\
& \leq \frac{1}{m}\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{2}\left(\sqrt{2}\|\widehat{D}-\bar{D}\|+\left\|\mathscr{P}_{\mathbb{T}}(-J(\widehat{D}-\bar{D}) J)\right\|_{*}+\left\|\mathscr{P}_{\mathbb{T}^{\perp}}(-J(\widehat{D}-\bar{D}) J)\right\|_{*}\right) \\
& \leq \frac{\rho_{D}}{\kappa_{D}}\left(\sqrt{2}\|\widehat{D}-\bar{D}\|+\left\|\mathscr{P}_{\mathbb{T}}(-J(\widehat{D}-\bar{D}) J)\right\|_{*}+\left\|\mathscr{P}_{\mathbb{T}^{\perp}}(-J(\widehat{D}-\bar{D}) J)\right\|_{*}\right) \\
& \leq \frac{\rho_{D}}{\kappa_{D}}\left(\sqrt{2}\|\widehat{D}-\bar{D}\|+\left\|\mathscr{P}_{\mathbb{T}}(-J(\widehat{D}-\bar{D}) J)\right\|_{*}+\sqrt{2 r}\|J(\widehat{D}-\bar{D}) J\|\right) \\
& \leq \frac{\rho_{D}}{\kappa_{D}}\left((\sqrt{2}+\sqrt{2 r})\|\widehat{D}-\bar{D}\|+\left\|\mathscr{P}_{\mathbb{T}}(-J(\widehat{D}-\bar{D}) J)\right\|_{*}\right)
\end{aligned}
$$

Similarly, we know from (23) and (25) that

$$
\frac{1}{m}\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{\infty}\|\widehat{S}-\bar{S}\|_{1} \leq \frac{\rho_{S}}{\kappa_{S}}\left(\left\|\mathscr{P}_{\mathbb{L}}(\widehat{S}-\bar{S})\right\|_{1}+\left\|\mathscr{P}_{\mathbb{L}^{\perp}}(\widehat{S}-\bar{S})\right\|_{1}\right) \leq \frac{\rho_{S}}{\kappa_{S}}\left(\sqrt{k}\|\widehat{S}-\bar{S}\|+\left\|\mathscr{P}_{\mathbb{L}^{\perp}}(\widehat{S}-\bar{S})\right\|_{1}\right)
$$

Therefore, we obtain that the first term of the right hand side of 29) satisfies

$$
\begin{align*}
\frac{1}{m}\left\langle\mathscr{O}_{\Omega}^{*}(\zeta), \widehat{D}-\bar{D}+\widehat{S}-\bar{S}\right\rangle \leq & \frac{\rho_{D}}{\kappa_{D}}\left(\sqrt{2}\|\widehat{D}-\bar{D}\|+\left\|\mathscr{P}_{\mathbb{T}}(-J(\widehat{D}-\bar{D}) J)\right\|_{*}+\left\|\mathscr{P}_{\mathbb{T}^{\perp}}(-J(\widehat{D}-\bar{D}) J)\right\|_{*}\right) \\
& +\frac{\rho_{S}}{\kappa_{S}}\left(\sqrt{k}\|\widehat{S}-\bar{S}\|+\left\|\mathscr{P}_{\mathbb{L}^{\perp}}(\widehat{S}-\bar{S})\right\|_{1}\right) \tag{30}
\end{align*}
$$

Meanwhile, since for any $A \in \mathbb{S}^{n},\left\|\mathscr{P}_{\mathbb{T}}(A)\right\|_{*}=\left\|\bar{P}_{2}^{T} A \bar{P}_{2}\right\|_{*}$ and both $-J \widetilde{D} J$ and $-J \bar{D} J$ are positively semidefinite, we know from the directional derivative formula of the nuclear norm [43, Theorem 1] that

$$
\begin{aligned}
\langle I,-J(\widehat{D}-\bar{D}) J\rangle=\|-J \widehat{D} J\|_{*}-\|-J \bar{D} J\|_{*} & \geq\left\langle\bar{P}_{1} \bar{P}_{1}^{T},-J(\widetilde{D}-\bar{D}) J\right\rangle+\left\|\bar{P}_{2}^{T}(-J(\widetilde{D}-\bar{D}) J) \bar{P}_{2}\right\|_{*} \\
& =\left\langle\bar{P}_{1} \bar{P}_{1}^{T},-J(\widetilde{D}-\bar{D}) J\right\rangle+\left\|\mathscr{P}_{\mathbb{T}}(-J(\widetilde{D}-\bar{D}) J)\right\|_{*}
\end{aligned}
$$

which implies that the second term of the right hand side of 29) satisfies

$$
-\rho_{D}(\langle I,-J(\widehat{D}-\bar{D}) J\rangle-\langle\widetilde{F},-J(\widehat{D}-\bar{D}) J\rangle) \leq-\rho_{D}\left(\left\langle\bar{P}_{1} \bar{P}_{1}^{T}-\widetilde{F},-J(\widetilde{D}-\bar{D}) J\right\rangle+\left\|\mathscr{P}_{\mathbb{T}}(-J(\widetilde{D}-\bar{D}) J)\right\|_{*}\right)
$$

By using the decomposition 20) and the notations defined in 24, we obtain that

$$
\begin{align*}
-\rho_{D}(\langle I,-J(\widehat{D}-\bar{D}) J\rangle-\langle\widetilde{F},-J(\widehat{D}-\bar{D}) J\rangle) & \leq \rho_{D}\left(\left\|\bar{P}_{1} \bar{P}_{1}^{T}-\widetilde{F}\right\|\|J(\widehat{D}-\bar{D}) J\|-\left\|\mathscr{P}_{\mathbb{T}}(-J(\widehat{D}-\bar{D}) J)\right\|_{*}\right) \\
& \leq \rho_{D}\left(a_{D} \sqrt{r}\|\widehat{D}-\bar{D}\|-\left\|\mathscr{P}_{\mathbb{T}}(-J(\widehat{D}-\bar{D}) J)\right\|_{*}\right) \tag{31}
\end{align*}
$$

Similarly, we know from the directional derivative of the $l_{1}$-norm at $\bar{S}$ that

$$
\|\widehat{S}\|_{1}-\|\bar{S}\|_{1} \geq\langle\operatorname{sign}(\bar{S}), \widehat{S}-\bar{S}\rangle+\left\|\mathscr{P}_{\mathbb{L}^{\perp}}(\widehat{S}-\bar{S})\right\|_{1}
$$

Therefore, by the decomposition (22) and the notations defined in 24, we know that the third term of the right hand side of 29) satisfies

$$
\begin{align*}
-\rho_{S}\left(\|\widehat{S}\|_{1}-\|\bar{S}\|_{1}-\langle\widetilde{G}, \widehat{S}-\bar{S}\rangle\right) & \leq-\rho_{S}\left(\langle\operatorname{sign}(\bar{S}), \widehat{S}-\bar{S}\rangle+\left\|\mathscr{P}_{\mathbb{L}^{\perp}}(\widehat{S}-\bar{S})\right\|_{1}-\langle\widetilde{G}, \widehat{S}-\bar{S}\rangle\right) \\
& \leq \rho_{S}\left(\|\operatorname{sign}(\bar{S})-\widetilde{G}\|\|\widehat{S}-\bar{S}\|-\left\|\mathscr{P}_{\mathbb{L}^{\perp}}(\widehat{S}-\bar{S})\right\|_{1}\right) \\
& \leq \rho_{S}\left(a_{S} \sqrt{k}\|\widehat{S}-\bar{S}\|-\left\|\mathscr{P}_{\mathbb{L}^{\perp}}(\widehat{S}-\bar{S})\right\|_{1}\right) \tag{32}
\end{align*}
$$

Finally, by substituting (30), (31) and (32) into (29), we obtain that

$$
\begin{align*}
& \frac{1}{2 m}\left\|\mathscr{O}_{\Omega}(\widehat{D}-\bar{D})+\mathscr{O}_{\Omega}(\widehat{S}-\bar{S})\right\|^{2} \\
\leq & \frac{\rho_{D}}{\kappa_{D}}\left((\sqrt{2}+\sqrt{2 r})\|\widehat{D}-\bar{D}\|+\left\|\mathscr{P}_{\mathbb{T}}(-J(\widehat{D}-\bar{D}) J)\right\|_{*}\right)+\frac{\rho_{S}}{\kappa_{S}}\left(\sqrt{k}\|\widehat{S}-\bar{S}\|+\left\|\mathscr{P}_{\mathbb{L}^{\perp}}(\widehat{S}-\bar{S})\right\|_{1}\right) \\
& +\rho_{D}\left(a_{D} \sqrt{r}\|\widehat{D}-\bar{D}\|-\left\|\mathscr{P}_{\mathbb{T}}(-J(\widehat{D}-\bar{D}) J)\right\|_{*}\right)+\rho_{S}\left(a_{S} \sqrt{k}\|\widehat{S}-\bar{S}\|-\left\|\mathscr{P}_{\mathbb{L}^{\perp}}(\widehat{S}-\bar{S})\right\|_{1}\right) \\
= & \rho_{D}\left(\frac{1}{\kappa_{D}}(\sqrt{2}+\sqrt{2 r})+a_{D} \sqrt{r}\right)\|\widehat{D}-\bar{D}\|+\rho_{S}\left(\frac{1}{\kappa_{S}}+a_{S}\right) \sqrt{k}\|\widehat{S}-\bar{S}\| \\
& -\rho_{D} \frac{\kappa_{D}-1}{\kappa_{D}}\left\|\mathscr{P}_{\mathbb{T}}(-J(\widehat{D}-\bar{D}) J)\right\|_{*}-\rho_{S} \frac{\kappa_{S}-1}{\kappa_{S}}\left\|\mathscr{P}_{\mathbb{L}^{\perp}}(\widehat{S}-\bar{S})\right\|_{1} . \tag{33}
\end{align*}
$$

Since $r \geq 1$, together with the assumptions that $\kappa_{D}>1$ and $\kappa_{S}>1$, we know that the inequality (26) holds.
Next, we shall show that the inequalities (27) also hold. By (33), we have

$$
\left\{\begin{array}{l}
\left\|\mathscr{P}_{\mathbb{T}}(-J(\widehat{D}-\bar{D}) J)\right\|_{*} \leq \frac{\kappa_{D}}{\kappa_{D}-1}\left(\sqrt{r}\left(a_{D}+\frac{2 \sqrt{2}}{\kappa_{D}}\right)\|\widehat{D}-\bar{D}\|+\sqrt{k} \frac{\rho_{S}}{\rho_{D}}\left(a_{S}+\frac{1}{\kappa_{S}}\right)\|\widehat{S}-\bar{S}\|\right) \\
\left\|\mathscr{P}_{\mathbb{L}^{\perp}}(\widehat{S}-\bar{S})\right\|_{1} \leq \frac{\kappa_{S}}{\kappa_{S}-1}\left(\sqrt{r} \frac{\rho_{D}}{\rho_{S}}\left(a_{D}+\frac{2 \sqrt{2}}{\kappa_{D}}\right)\|\widehat{D}-\bar{D}\|+\sqrt{k}\left(a_{S}+\frac{1}{\kappa_{S}}\right)\|\widehat{S}-\bar{S}\|\right)
\end{array}\right.
$$

Therefore, we know from 20, 22, 22, and 23) that

$$
\begin{aligned}
\|\widehat{D}-\bar{D}\|_{*} & \leq\|\widehat{D}-\bar{D}-J(\widehat{D}-\bar{D}) J\|_{*}+\left\|\mathscr{P}_{\mathbb{T}^{\perp}}(-J(\widehat{D}-\bar{D}) J)\right\|_{*}+\left\|\mathscr{P}_{\mathbb{T}}(-J(\widehat{D}-\bar{D}) J)\right\|_{*} \\
& \leq(\sqrt{2}+\sqrt{2 r})\|\widehat{D}-\bar{D}\|+\frac{\kappa_{D}}{\kappa_{D}-1}\left(\sqrt{r}\left(a_{D}+\frac{2 \sqrt{2}}{\kappa_{D}}\right)\|\widehat{D}-\bar{D}\|+\sqrt{k} \frac{\rho_{S}}{\rho_{D}}\left(a_{S}+\frac{1}{\kappa_{S}}\right)\|\widehat{S}-\bar{S}\|\right) \\
& \leq \frac{\kappa_{D}}{\kappa_{D}-1}\left(\left(a_{D}+2 \sqrt{2}\right) \sqrt{r}\|\widehat{D}-\bar{D}\|+\frac{\rho_{S}}{\rho_{D}}\left(a_{S}+\frac{1}{\kappa_{S}}\right) \sqrt{k}\|\widehat{S}-\bar{S}\|\right)
\end{aligned}
$$

and

$$
\begin{aligned}
\|\widehat{S}-\bar{S}\|_{1} & \leq\left\|\mathscr{P}_{\mathbb{L}}(\widehat{S}-\bar{S})\right\|_{1}+\left\|\mathscr{P}_{\mathbb{L}^{\perp}}(\widehat{S}-\bar{S})\right\|_{1} \\
& \leq \sqrt{k}\|\widehat{S}-\bar{S}\|+\frac{\kappa_{S}}{\kappa_{S}-1}\left(\sqrt{r} \frac{\rho_{D}}{\rho_{S}}\left(a_{D}+\frac{2 \sqrt{2}}{\kappa_{D}}\right)\|\widehat{D}-\bar{D}\|+\sqrt{k}\left(a_{S}+\frac{1}{\kappa_{S}}\right)\|\widehat{S}-\bar{S}\|\right) \\
& \leq \frac{\kappa_{S}}{\kappa_{S}-1}\left(\frac{\rho_{D}}{\rho_{S}}\left(a_{D}+\frac{2 \sqrt{2}}{\kappa_{D}}\right) \sqrt{r}\|\widehat{D}-\bar{D}\|+\left(a_{S}+1\right) \sqrt{k}\|\widehat{S}-\bar{S}\|\right) .
\end{aligned}
$$

This completes the proof.
Since $X_{1}, \ldots, X_{m}$ are the i.i.d. random observations, i.e., for any $1 \leq i<j \leq n$,

$$
\mathbb{P}\left(X_{l}=\frac{1}{2}\left(\mathbf{e}_{i} \mathbf{e}_{j}^{T}+\mathbf{e}_{j} \mathbf{e}_{i}^{T}\right)\right)=\pi_{i j}, \quad l=1, \ldots, m
$$

where $0 \leq \pi_{i j} \leq 1$ is the probability that the $(i, j)$ and $(j, i)$-th element be sampled in the observation model. We propose the following assumption to control the sampling probability.

Assumption 2 There exist two absolution constants $\mu_{1}, \mu_{2} \geq 1$ such that

$$
\frac{1}{\mu_{1} d_{\mathbb{S}_{h}^{n}}} \leq \pi_{i j} \leq \frac{\mu_{2}}{d_{\mathbb{S}_{h}^{n}}} \quad \forall 1 \leq i<j \leq n
$$

where $d_{\mathbb{S}_{h}^{n}}=n(n-1) / 2$.
It is easy to see from Assumption 2 that for any $A \in \mathbb{S}_{h}^{n}$, we have

$$
\begin{equation*}
\mathbb{E}\left(\langle A, X\rangle^{2}\right) \geq \frac{1}{2 \mu_{1} d_{\mathbb{S}_{h}^{n}}}\|A\|^{2} \tag{34}
\end{equation*}
$$

Furthermore, let $m_{\max }$ be the maximum number of repetitions of any $(i, j)$ index in $\Omega$. By noting the sample size $m$ is assumed much smaller than $d_{\mathbb{S}_{h}^{n}}$, we obtain from [46, Lemma 5.5] the following result on the upper bound of $m_{\text {max }}$. For simplicity, we omit the detailed proof here.

Lemma 2 Let the observation index set $\Omega$ be generated by the uniform sampling with replacement. Then, there exists a constant $C>0$ such that

$$
m_{\max } \leq\left\|\mathscr{O}_{\Omega}^{*} \mathscr{O}_{\Omega}\right\|_{2} \leq C \log \left(2 n^{2}\right)
$$

with probability at least $1-1 /\left(2 n^{2}\right)$.

We further introduce the following two useful notations:

$$
\begin{equation*}
\vartheta_{D}:=\mathbb{E}\left\|\frac{1}{m} \mathscr{O}_{\Omega}^{*}(\varepsilon)\right\|_{2} \quad \text { and } \quad \vartheta_{S}:=\mathbb{E}\left\|\frac{1}{m} \mathscr{O}_{\Omega}^{*}(\varepsilon)\right\|_{\infty} \tag{35}
\end{equation*}
$$

where $\left\{\varepsilon_{1}, \ldots, \varepsilon_{m}\right\}$ is a Rademacher sequence, i.e., an i.i.d. copy of Bernoulli random variable taking the values 1 and -1 with probability $1 / 2$.

For the given positive numbers $p_{1}, p_{2}, q_{1}, q_{2}$ and $t$, define the following subset $K\left(p_{1}, p_{2}, q_{1}, q_{2}, t\right) \subseteq \mathbb{S}^{n}$ by

$$
K(\mathbf{p}, \mathbf{q}, t):=\left\{\begin{array}{ll} 
& \left\|A_{D}\right\|_{*} \leq p_{1}\left\|A_{D}\right\|+p_{2}\left\|A_{S}\right\|, A_{D} \in \mathbb{S}_{h}^{n}  \tag{36}\\
A=A_{D}+A_{S} \mid & \left\|A_{S}\right\|_{1} \leq q_{1}\left\|A_{D}\right\|+q_{2}\left\|A_{S}\right\|, A_{S} \in \mathbb{S}_{h}^{n} \\
& \left\|A_{D}+A_{S}\right\|_{\infty}=1,\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2} \geq t \mu_{1} d_{\mathbb{S}_{h}^{n}}
\end{array}\right\}
$$

where $\mathbf{p}:=\left(p_{1}, p_{2}\right)$ and $\mathbf{q}:=\left(q_{1}, q_{2}\right)$. Denote $\vartheta_{m}:=\left(\vartheta_{D}^{2} p_{1}^{2}+\vartheta_{D}^{2} p_{2}^{2}+\vartheta_{S}^{2} p_{1}^{2}+\vartheta_{S}^{2} p_{2}^{2}\right)$.
Proposition 2 Suppose that Assumption 2 holds. Let $p_{1}, p_{2}, q_{1}, q_{2}$ and $t$ be any given positive numbers. For any $\tau_{1}$ and $\tau_{2}$ satisfying

$$
0<\tau_{1}<1 \quad \text { and } \quad 0<\tau_{2}<\frac{\tau_{1}}{2}
$$

it holds that for any $A \in K(\mathbf{p}, \mathbf{q}, t)$,

$$
\begin{equation*}
\frac{1}{m}\left\|\mathscr{O}_{\Omega}(A)\right\|^{2} \geq \mathbb{E}\left(\langle A, X\rangle^{2}\right)-\frac{\tau_{1}}{\mu_{1} d_{\mathbb{S}_{h}^{n}}}\left(\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}\right)-\frac{32}{\tau_{2}} \mu_{1} d_{\mathbb{S}_{h}^{n}} \vartheta_{m}^{2} \tag{37}
\end{equation*}
$$

with probability at least

$$
1-\frac{\exp \left(-\left(\tau_{1}-2 \tau_{2}\right)^{2} m t^{2} / 8\right)}{1-\exp \left(-3\left(\tau_{1}-2 \tau_{2}\right)^{2} m t^{2} / 8\right)}
$$

Proof We will show that the event

$$
E:=\left\{\exists A \in K(\mathbf{p}, \mathbf{q}, t) \text { s.t. }\left|\frac{1}{m}\left\|\mathscr{O}_{\Omega}(A)\right\|^{2}-\mathbb{E}\left(\langle A, X\rangle^{2}\right)\right| \geq \frac{\tau_{1}}{\mu_{1} d_{\mathbb{S}_{h}^{n}}}\left(\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}\right)+\frac{32}{\tau_{2}} \mu_{1} d_{\mathbb{S}_{h}^{n}} \vartheta_{m}^{2}\right\}
$$

happens with probability less than $\frac{\exp \left[-\left(\tau_{1}-2 \tau_{2}\right)^{2} m t^{2} / 32\right]}{1-\exp \left[-3\left(\tau_{1}-2 \tau_{2}\right)^{2} m t^{2} / 32\right]}$. First, we decompose $K(\mathbf{p}, \mathbf{q}, t)$ by

$$
K(\mathbf{p}, \mathbf{q}, t)=\bigcup_{j=1}^{\infty}\left\{A \in K(\mathbf{p}, \mathbf{q}, t) \left\lvert\, 2^{j-1} t \leq \frac{1}{\mu_{1} d_{\mathbb{S}_{h}^{n}}}\left(\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}\right) \leq 2^{j} t\right.\right\}
$$

For any $s \geq t$, define the sunset $\widetilde{K}(\mathbf{p}, \mathbf{q}, t, s) \subseteq K(\mathbf{p}, \mathbf{q}, t)$ by

$$
\widetilde{K}(\mathbf{p}, \mathbf{q}, t, s):=\left\{A \in K(\mathbf{p}, \mathbf{q}, t) \left\lvert\, \frac{1}{\mu_{1} d_{\mathbb{S}_{h}^{n}}}\left(\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}\right) \leq s\right.\right\} .
$$

Furthermore, for $j=1,2, \ldots$, let $E_{j}$ be the set defined by

$$
E_{j}:=\left\{\exists A \in \widetilde{K}\left(\mathbf{p}, \mathbf{q}, t, 2^{j} t\right) \text { such that }\left|\frac{1}{m}\left\|\mathscr{O}_{\Omega}(A)\right\|^{2}-\mathbb{E}\left(\langle A, X\rangle^{2}\right)\right| \geq \tau_{1} 2^{j-1} t+\frac{32}{\tau_{2}} \mu_{1} d_{\mathbb{S}_{h}^{n}} \vartheta_{m}^{2}\right\} .
$$

Then, it is not difficult to see that $E \subseteq \cup_{j=1}^{\infty} E_{j}$. Thus, it suffices to estimate the probability of each simpler event $E_{j}$ and then obtain the estimated probability bound of the event $E$. Denote

$$
Z_{s}:=\sup _{A \in \widetilde{K}(\mathbf{p}, \mathbf{q}, t, s)}\left|\frac{1}{m}\left\|\mathscr{O}_{\Omega}(A)\right\|^{2}-\mathbb{E}\left(\langle A, X\rangle^{2}\right)\right| .
$$

For any $A \in \mathbb{S}_{h}^{n}$, the strong laws of large numbers yield that

$$
\frac{1}{m}\left\|\mathscr{O}_{\Omega}(A)\right\|^{2}=\frac{1}{m} \sum_{l=1}^{m}\left\langle X_{l}, A\right\rangle^{2} \xrightarrow{\text { a.s. }} \mathbb{E}\left(\langle A, X\rangle^{2}\right) \quad \text { as } \quad m \rightarrow \infty .
$$

Since $\|A\|_{\infty}=1$ for all $A \in K(\mathbf{p}, \mathbf{q}, t)$, we know that for any $1 \leq l \leq m$ and $A \in K(\mathbf{p}, \mathbf{q}, t)$,

$$
\left|\left\langle X_{l}, A\right\rangle^{2}-\mathbb{E}\left(\left\langle X_{l}, A\right\rangle^{2}\right)\right| \leq \max \left\{\left\langle X_{l}, A\right\rangle^{2}, \mathbb{E}\left(\left\langle X_{l}, A\right\rangle^{2}\right)\right\} \leq 1
$$

Then, according to Massart's Hoeffding-type concentration inequality [2, Theorem 14.2] (see also [26, Theorem 9]), we know that

$$
\begin{equation*}
\mathbb{P}\left(Z_{s} \leq \mathbb{E}\left(Z_{s}\right)+\eta\right) \leq \exp \left(-\frac{m \eta^{2}}{8}\right) \quad \forall \eta \geq 0 \tag{38}
\end{equation*}
$$

Next, we estimate an upper bound of $\mathbb{E}\left(Z_{S}\right)$ by using the standard Rademacher symmetrization in the theory of empirical processes. Recall that $\left\{\varepsilon_{1}, \ldots, \varepsilon_{m}\right\}$ is a Rademacher sequence. Then, we have

$$
\begin{aligned}
\mathbb{E}\left(Z_{s}\right) & =\mathbb{E}\left(\sup _{A \in \widetilde{K}(\mathbf{p}, \mathbf{q}, t, s)}\left|\frac{1}{m} \sum_{l=1}^{m}\left\langle X_{l}, A\right\rangle^{2}-\mathbb{E}\left[\left\langle X_{l}, A\right\rangle^{2}\right]\right|\right) \leq 2 \mathbb{E}\left(\sup _{A \in \widetilde{K}(\mathbf{p}, \mathbf{q}, t, s)}\left|\frac{1}{m} \sum_{l=1}^{m} \varepsilon_{l}\left\langle X_{l}, A\right\rangle^{2}\right|\right) \\
& \leq 8 \mathbb{E}\left(\sup _{A \in \widetilde{K}(\mathbf{p}, \mathbf{q}, t, s)}\left|\frac{1}{m} \sum_{l=1}^{m} \varepsilon_{l}\left\langle X_{l}, A\right\rangle\right|\right)=8 \mathbb{E}\left(\sup _{A \in \widetilde{K}(\mathbf{p}, \mathbf{q}, t, s)}\left|\left\langle\frac{1}{m} \mathscr{O}^{*}(\varepsilon), A\right\rangle\right|\right) \\
& \leq 8 \mathbb{E}\left(\sup _{A \in \widetilde{K}(\mathbf{p}, \mathbf{q}, t, s)}\left(\left\|\frac{1}{m} \mathscr{O}^{*}(\varepsilon)\right\|_{2}\left\|A_{D}\right\|_{*}+\left\|\frac{1}{m} \mathscr{O}^{*}(\varepsilon)\right\|_{\infty}\left\|A_{S}\right\|_{1}\right)\right) \\
& \leq 8 \mathbb{E}\left(\sup _{A \in \widetilde{K}(\mathbf{p}, \mathbf{q}, t, s)}\left\|\frac{1}{m} \mathscr{O}^{*}(\varepsilon)\right\|_{2}\left\|A_{D}\right\|_{*}+\sup _{A \in \widetilde{K}(\mathbf{p}, \mathbf{q}, t, s)}\left\|\frac{1}{m} \mathscr{O}^{*}(\varepsilon)\right\|_{\infty}\left\|A_{S}\right\|_{1}\right) \\
& \leq 8 \mathbb{E}\left\|\frac{1}{m} \mathscr{O}^{*}(\varepsilon)\right\|_{2}\left(\sup _{A \in \widetilde{K}(\mathbf{p}, \mathbf{q}, t, s)}\left\|A_{D}\right\|_{*}\right)+8 \mathbb{E}\left\|\frac{1}{m} \mathscr{O}^{*}(\varepsilon)\right\|_{\infty}\left(\sup _{A \in \widetilde{K}(\mathbf{p}, \mathbf{q}, t, s)}\left\|A_{S}\right\|_{1}\right),
\end{aligned}
$$

where the first inequality is due to the symmetrization theorem [39] Lemma 2.3.1] or [2, Theorem 14.3], and the second inequality follows from the contraction theorem (e.g., [24, Theorem 4.12] and [2, Theorem 14.4], ). Notice that for any $u \geq 0, v \geq 0$ and $A \in \widetilde{K}(\mathbf{p}, \mathbf{q}, t, s)$,

$$
u\left\|A_{D}\right\|+v\left\|A_{S}\right\| \leq \frac{16 m}{\tau_{1}}\left(u^{2}+v^{2}\right)+\frac{\tau_{1}}{64 m}\|A\|^{2} \leq \frac{16 m}{\tau_{1}}\left(u^{2}+v^{2}\right)+\frac{1}{64} \tau_{1} s
$$

where the first inequality is due to the inequality of arithmetic and geometric means. We derive that

$$
\begin{aligned}
\mathbb{E}\left(Z_{S}\right) & \leq 8\left(\sup _{A \in \widetilde{K}(\mathbf{p}, \mathbf{q}, t, s)} \vartheta_{D}\left(p_{1}\left\|A_{D}\right\|+p_{2}\left\|A_{S}\right\|\right)+\sup _{A \in \widetilde{K}(\mathbf{p}, \mathbf{q}, t, s)} \vartheta_{S}\left(q_{1}\left\|A_{D}\right\|+q_{2}\left\|A_{S}\right\|\right)\right) \\
& \leq \frac{16}{\tau_{1}} m\left(\vartheta_{D}^{2} p_{1}^{2}+\vartheta_{D}^{2} p_{2}^{2}+\vartheta_{S}^{2} p_{1}^{2}+\vartheta_{S}^{2} p_{2}^{2}\right)+\frac{\tau_{1}}{32} s=\frac{16}{\tau_{1}} m \vartheta_{m}^{2}+\frac{\tau_{1}}{32} s .
\end{aligned}
$$

According to 38, it follows that

$$
\mathbb{P}\left(Z_{s} \geq \frac{16}{\tau_{1}} m \vartheta_{m}^{2}+\frac{\tau_{1}}{32} s\right) \leq \mathbb{P}\left(Z_{s} \geq \mathbb{E}\left(Z_{s}\right)+\left(\frac{\tau_{1}}{2}-\tau_{2}\right) s\right) \leq \exp \left(-\left(\frac{\tau_{1}}{2}-\tau_{2}\right)^{2} \frac{m s^{2}}{8}\right)
$$

This, together with the choice of $s=2^{j} t$, implies that $\mathbb{P}\left(E_{j}\right) \leq \exp \left(-\frac{1}{8} 2^{2(j-1)}\left(\tau_{1}-2 \tau_{2}\right)^{2} m t^{2}\right)$. By using the fact that $2^{j} \geq 1+j(2-1)$ for any $j \geq 1$, we obtain that

$$
\begin{aligned}
\mathbb{P}(E) & \leq \sum_{j=1}^{\infty} \mathbb{P}\left(E_{j}\right) \leq \sum_{j=1}^{\infty} \exp \left(-\frac{1}{8} 2^{2(j-1)}\left(\tau_{1}-2 \tau_{2}\right)^{2} m t^{2}\right) \\
& \leq \exp \left(-\frac{1}{8}\left(\tau_{1}-2 \tau_{2}\right)^{2} m t^{2}\right) \sum_{j=1}^{\infty} \exp \left(-\frac{1}{8}\left(2^{2(j-1)}-1\right)\left(\tau_{1}-2 \tau_{2}\right)^{2} m t^{2}\right) \\
& \leq \exp \left(-\frac{1}{8}\left(\tau_{1}-2 \tau_{2}\right)^{2} m t^{2}\right) \sum_{j=1}^{\infty} \exp \left(-\frac{3}{8}(j-1)\left(\tau_{1}-2 \tau_{2}\right)^{2} m t^{2}\right) . \\
& =\frac{\exp \left(-\left(\tau_{1}-2 \tau_{2}\right)^{2} m t^{2} / 8\right)}{1-\exp \left(-3\left(\tau_{1}-2 \tau_{2}\right)^{2} m t^{2} / 8\right)} .
\end{aligned}
$$

The proof is then completed.
Proposition 3 Let $(\widehat{D}, \widehat{S})$ and $(\bar{D}, \bar{S})$ be an optimal solution of $\sqrt{12}$ ) and the underground true EDM and outlier matrices, respectively. Let $\kappa_{D}>1$ and $\kappa_{S}>1$ be given arbitrarily. Suppose that the parameters $\rho_{D}>0$ and $\rho_{S}>0$ are given by (25). Under Assumption 2, there exist some positive absolute constants $C_{0}, C_{1}$ and $C_{2}$ such that either

$$
\begin{equation*}
\frac{\|\widehat{D}-\bar{D}\|^{2}+\|\widehat{S}-\bar{S}\|^{2}}{d_{\mathbb{S}_{h}^{n}}} \leq C_{0} \mu_{1}\left(b_{D}+b_{S}\right)^{2} \sqrt{\frac{\log (2 n)}{m}} \tag{39}
\end{equation*}
$$

or

$$
\begin{aligned}
\frac{\|\widehat{D}-\bar{D}\|^{2}+\|\widehat{S}-\bar{S}\|^{2}}{d_{\mathbb{S}_{h}^{n}} \leq} & C_{1} \mu_{1}^{2} d_{\mathbb{S}_{h}^{n}}\left\{C_{2}^{2}\left[\rho_{D}^{2} r\left(a_{D}+\frac{2 \sqrt{2}}{\kappa_{D}}\right)^{2}+\rho_{S}^{2} k\left(a_{S}+\frac{1}{\kappa_{S}}\right)^{2}\right]\right. \\
& +\vartheta_{D}^{2}\left(b_{D}+b_{S}\right)^{2}\left(\frac{\kappa_{D}}{\kappa_{D}-1}\right)^{2}\left[r\left(a_{D}+2 \sqrt{2}\right)^{2}+k \frac{\rho_{S}^{2}}{\rho_{D}^{2}}\left(a_{S}+\frac{1}{\kappa_{S}}\right)^{2}\right] \\
& \left.+\max \left\{\vartheta_{S}^{2}\left(b_{D}+b_{S}\right)^{2}, \frac{b_{D}^{2}}{\mu_{1}^{2} d_{\mathbb{S}_{h}^{2}}^{2}}\right\}\left(\frac{\kappa_{S}}{\kappa_{S}-1}\right)^{2}\left[r \frac{\rho_{D}^{2}}{\rho_{S}^{2}}\left(a_{D}+\frac{2 \sqrt{2}}{\kappa_{D}}\right)^{2}+k\left(a_{S}+1\right)^{2}\right]\right\}(40)
\end{aligned}
$$

with probability at least $1-(4 / 7) n^{-1}$, where $a_{D}$ and $a_{S}$ are given by (24) and $\vartheta_{D}$, and $\vartheta_{S}$ are defined by (35).
Proof Denote $A:=A_{D}+A_{S}$ with $A_{D}:=\widehat{D}-\bar{D}$ and $A_{S}:=\widehat{S}-\bar{S}$. Let $b:=\|A\|_{\infty}$ and $t:=\sqrt{\frac{32 \log (2 n)}{\left(\tau_{1}-2 \tau_{2}\right)^{2} m}}$, where $c>0$ and $\tau_{1}, \tau_{2}$ satisfying $0<\tau_{2}<\frac{\tau_{1}}{2}<1 / 2$ are arbitrarily fixed constants. Consider the following two cases.

Case 1. $\left(\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}\right)<b^{2} \mu_{1} d_{\mathbb{S}_{h}^{n} t}$. Since $b=\|A\|_{\infty} \leq\left\|A_{D}\right\|_{\infty}+\left\|A_{S}\right\|_{\infty} \leq 2\left(b_{D}+b_{S}\right)$, we know that there exists a positive constant $C_{0}$ such that 39 holds.

Case 2. $\left(\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}\right) \geq b^{2} \mu_{1} d_{\mathbb{S}_{h}^{n}}$. By 27) in Proposition 1 , we know that $A / b \in K(\mathbf{p}, \mathbf{q}, t)$, where $K(\mathbf{p}, \mathbf{q}, t)$ is the subset defined by with $\mathbf{p}=\left(p_{1}, p_{2}\right)$ and $\mathbf{q}=\left(q_{1}, q_{2}\right)$ are given by

$$
\begin{cases}p_{1}=\frac{\kappa_{D}}{\kappa_{D}-1}\left(a_{D}+2 \sqrt{2}\right) \sqrt{r}, & p_{2}=\frac{\kappa_{D}}{\kappa_{D}-1} \frac{\rho_{S}}{\rho_{D}}\left(a_{S}+\frac{1}{\kappa_{S}}\right) \sqrt{k}  \tag{41}\\ q_{1}=\frac{\kappa_{S}}{\kappa_{S}-1} \frac{\rho_{D}}{\rho_{S}}\left(a_{D}+\frac{2 \sqrt{2}}{\kappa_{D}}\right) \sqrt{r}, & q_{2}=\frac{\kappa_{S}}{\kappa_{S}-1}\left(a_{S}+1\right) \sqrt{k}\end{cases}
$$

Therefore, it follows from Proposition 2 and $\sqrt{34}$ that with probability at least $1-(4 / 7) n^{-1}$,

$$
\frac{1}{d_{\mathbb{S}_{h}^{n}}}\|A\|^{2} \leq 2 \mu_{1} \mathbb{E}\left(\langle A, X\rangle^{2}\right) \leq \frac{2 \mu_{1}}{m}\left\|\mathscr{O}_{\Omega}(A)\right\|^{2}+\frac{2 \tau_{1}}{d_{\mathbb{S}_{h}^{n}}}\left(\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}\right)+\frac{64}{\tau_{2}} \mu_{1}^{2} d_{\mathbb{S}_{h}^{n}} \vartheta_{m}^{2} b^{2}
$$

By (26) in Proposition 1, we obtain that for any $0<\tau_{3}<\left(1-2 \tau_{1}\right) / 2$,

$$
\begin{align*}
\frac{1}{d_{\mathbb{S}_{h}^{n}}}\|A\|^{2} \leq & 4 \mu_{1} \rho_{D} \sqrt{r}\left(a_{D}+\frac{2 \sqrt{2}}{\kappa_{D}}\right)\left\|A_{D}\right\|+4 \mu_{1} \rho_{S} \sqrt{k}\left(a_{S}+\frac{1}{\kappa_{S}}\right)\left\|A_{S}\right\| \\
& +\frac{2 \tau_{1}}{d_{\mathbb{S}_{h}^{n}}}\left(\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}\right)+\frac{64}{\tau_{2}} \mu_{1}^{2} d_{\mathbb{S}_{h}^{n}} \vartheta_{m}^{2} b^{2} \\
\leq & \frac{4 \mu_{1}^{2} \rho_{D}^{2} r d_{\mathbb{S}_{h}^{n}}}{\tau_{3}}\left(a_{D}+\frac{2 \sqrt{2}}{\kappa_{D}}\right)^{2}+\frac{\tau_{3}}{d_{\mathbb{S}_{h}^{n}}}\left\|A_{D}\right\|^{2}+\frac{4 \mu_{1}^{2} \rho_{S}^{2} k d_{\mathbb{S}_{h}^{n}}}{\tau_{3}}\left(a_{S}+\frac{1}{\kappa_{S}}\right)^{2}+\frac{\tau_{3}}{d_{\mathbb{S}_{h}^{n}}}\left\|A_{S}\right\|^{2} \\
& +\frac{2 \tau_{1}}{d_{\mathbb{S}_{h}^{n}}}\left(\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}\right)+\frac{64}{\tau_{2}} \mu_{1}^{2} d_{\mathbb{S}_{h}^{n}} \vartheta_{m}^{2} b^{2} \\
= & \frac{4 \mu_{1}^{2} \rho_{D}^{2} r d_{\mathbb{S}_{h}^{n}}}{\tau_{3}}\left(a_{D}+\frac{2 \sqrt{2}}{\kappa_{D}}\right)^{2}+\frac{4 \mu_{1}^{2} \rho_{S}^{2} k d_{\mathbb{S}_{h}^{n}}}{\tau_{3}}\left(a_{S}+\frac{1}{\kappa_{S}}\right)^{2}+\frac{2 \tau_{1}+\tau_{3}}{d_{\mathbb{S}_{h}^{n}}}\left(\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}\right) \\
& +\frac{64}{\tau_{2}} \mu_{1}^{2} d_{\mathbb{S}_{h}^{n}} \vartheta_{m}^{2} b^{2} . \tag{42}
\end{align*}
$$

In addition, since $\left\|A_{D}\right\| \leq 2 b_{D}$, we then derive from that

$$
\begin{aligned}
\|A\|^{2} & \geq\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}-2\left\|A_{D}\right\|_{\infty}\left\|A_{S}\right\|_{1} \geq\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}-4 b_{D}\left(q_{1}\left\|A_{D}\right\|+q_{2}\left\|A_{S}\right\|\right) \\
& \geq\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}-\frac{4}{\tau_{3}} b_{D}^{2}\left(q_{1}^{2}+q_{2}^{2}\right)-\tau_{3}\left(\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}\right)
\end{aligned}
$$

This, together with (42), yields that

$$
\begin{aligned}
\frac{1-\tau_{3}}{d_{\mathbb{S}_{h}^{n}}}\left(\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}\right) \leq & \frac{1}{d_{\mathbb{S}_{h}^{n}}}\|A\|^{2}+\frac{4}{d_{\mathbb{S}_{h}^{n}} \tau_{3}} b_{D}^{2}\left(q_{1}^{2}+q_{2}^{2}\right) \\
\leq & \frac{4 \mu_{1}^{2} \rho_{D}^{2} r d_{\mathbb{S}_{h}^{n}}}{\tau_{3}}\left(a_{D}+\frac{2 \sqrt{2}}{\kappa_{D}}\right)^{2}+\frac{4 \mu_{1}^{2} \rho_{S}^{2} k d_{\mathbb{S}_{h}^{n}}}{\tau_{3}}\left(a_{S}+\frac{1}{\kappa_{S}}\right)^{2}+\frac{2 \tau_{1}+\tau_{3}}{d_{\mathbb{S}_{h}^{n}}}\left(\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}\right) \\
& +\frac{64}{\tau_{2}} \mu_{1}^{2} d_{\mathbb{S}_{h}^{n}} \vartheta_{m}^{2} b^{2}+\frac{4}{d_{\mathbb{S}_{h}^{n}} \tau_{3}} b_{D}^{2}\left(q_{1}^{2}+q_{2}^{2}\right) \\
= & \frac{4 \mu_{1}^{2} d_{\mathbb{S}_{h}^{n}}}{\tau_{3}}\left(\rho_{D}^{2} r\left(a_{D}+\frac{2 \sqrt{2}}{\kappa_{D}}\right)^{2}+\rho_{S}^{2} k\left(a_{S}+\frac{1}{\kappa_{S}}\right)^{2}\right)+\frac{2 \tau_{1}+\tau_{3}}{d_{\mathbb{S}_{h}^{n}}}\left(\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}\right) \\
& +\frac{64}{\tau_{2}} \mu_{1}^{2} d_{\mathbb{S}_{h}^{n}} \vartheta_{m}^{2} b^{2}+\frac{4}{d_{\mathbb{S}_{h}^{n}} \tau_{3}} b_{D}^{2}\left(q_{1}^{2}+q_{2}^{2}\right) .
\end{aligned}
$$

Since $1-2\left(\tau_{1}+\tau_{3}\right)>0$, we have

$$
\begin{aligned}
\frac{\left\|A_{D}\right\|^{2}+\left\|A_{S}\right\|^{2}}{d_{\mathbb{S}_{h}^{n}}} \leq & \frac{4 \mu_{1}^{2} d_{\mathbb{S}_{h}^{n}}}{1-2\left(\tau_{1}+\tau_{3}\right)}\left(\frac{1}{\tau_{3}}\left(\rho_{D}^{2} r\left(a_{D}+\frac{2 \sqrt{2}}{\kappa_{D}}\right)^{2}+\rho_{S}^{2} k\left(a_{S}+\frac{1}{\kappa_{S}}\right)^{2}\right)+\frac{16}{\tau_{2}} \vartheta_{m}^{2} b^{2}\right. \\
& \left.+\frac{1}{d_{\mathbb{S}_{h}^{n}}^{2} \mu_{1}^{2} \tau_{3}} b_{D}^{2}\left(q_{1}^{2}+q_{2}^{2}\right)\right)
\end{aligned}
$$

Recall that $\vartheta_{m}^{2}=\vartheta_{D}^{2} p_{1}^{2}+\vartheta_{D}^{2} p_{2}^{2}+\vartheta_{S}^{2} q_{1}^{2}+\vartheta_{S}^{2} q_{2}^{2}$. By plugging this together with 41 into the above inequality and choosing $\tau_{1}, \tau_{2}$ and $\tau_{3}$ to be constants, we complete the proof.

In order to obtain the explicit formulas of the penalized parameters $\rho_{D}$ and $\rho_{S}$ based on 25], we shall derive the probabilistic upper bounds on the terms $\frac{1}{m}\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{2}$ and $\frac{1}{m}\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{\infty}$. To this end, similar with [11], from now on, we always assume that the i.i.d. random noises $\xi_{l}, l=1, \ldots, m$ in the sampling model (5) satisfy the following sub-Gaussian tail condition.

Assumption 3 There exist positive constants $K_{1}$ and $K_{2}$ such that for all $t>0$,

$$
\mathbb{P}\left(\left|\xi_{l}\right| \geq t\right) \leq K_{1} \exp \left(-t^{2} / K_{2}\right)
$$

The following proposition on the upper bounds on the terms $\frac{1}{m}\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{2}$ is taken from [11] Proposition 4].
Proposition 4 Let $\zeta=\mathscr{O}_{\Omega}(ד) \in \mathbb{R}^{m}$ and $\rceil \in \mathbb{S}^{n}$ be given by (8). Suppose that there exists $C_{1}>1$ such that $m>C_{1} n \log (n)$. Then, there exists a constant $C_{2}>0$ such that with probability at least $1-1 / n$,

$$
\begin{equation*}
\frac{1}{m}\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{2} \leq C_{2}\left(2 \omega \eta+\eta^{2}\right) \sqrt{\frac{\log (2 n)}{n m}} \tag{43}
\end{equation*}
$$

where $\omega=\left\|\mathscr{O}_{\Omega}(\bar{d}+\bar{s})\right\|_{\infty}$.
The following result on the upper bound of $\frac{1}{m}\left\|\mathscr{O}^{*}(\xi)\right\|_{\infty}$ are a direct consequence of the large derivation inequality for sums of independent sub-gaussian/sub-exponential random variables [42, Proposition 5.10 \& 5.16].

Proposition 5 Let $\zeta=\mathscr{O}_{\Omega}(ד) \in \mathbb{R}^{m}$ and $\rceil \in \mathbb{S}^{n}$ be given by (8). Then, there exists a positive constant $C_{3}$ such that with probability at least $1-2 / n^{2}$,

$$
\begin{equation*}
\frac{1}{m}\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{\infty} \leq C_{3}\left(2 \omega \eta+\eta^{2}\right) \frac{\log \left(2 n^{2}\right)}{m} \tag{44}
\end{equation*}
$$

where $\omega=\left\|\mathscr{O}_{\Omega}(\bar{d}+\bar{s})\right\|_{\infty}$.
Proof From 8 and the definition of $\zeta$, we know that

$$
\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{\infty} \leq 2 \omega \eta\left\|\mathscr{O}_{\Omega}^{*}(\xi)\right\|_{\infty}+\eta^{2}\left\|\mathscr{O}_{\Omega}^{*}(\xi \circ \xi)\right\|_{\infty} .
$$

Therefore, for any given $t_{1}, t_{2}>0$, we have

$$
\begin{equation*}
\mathbb{P}\left(\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{\infty} \geq 2 \omega \eta t_{1}+\eta^{2} t_{2}\right) \leq \mathbb{P}\left(\left\|\mathscr{O}_{\Omega}^{*}(\xi)\right\|_{\infty} \geq t_{1}\right)+\mathbb{P}\left(\left\|\mathscr{O}_{\Omega}^{*}(\xi \circ \xi)\right\|_{\infty} \geq t_{2}\right) \tag{45}
\end{equation*}
$$

Denote the random matrix $Y:=\mathscr{O}_{\Omega}^{*}(\xi)=\sum_{l=1}^{m} \xi_{l} X_{l}$ and $Z:=\mathscr{O}_{\Omega}^{*}(\xi \circ \xi-\mathbf{1})=\sum_{l=1}^{m}\left(\xi_{l}^{2}-1\right) X_{l}$. Then, for each $i, j \in\{1, \ldots, n\}$, the $(i, j)$-th elements of $Y$ and $Z$ can be written as $Y_{i j}=\sum_{l=1}^{m} \mathbf{a}_{l}^{(i j)} \xi_{l}$ and $Z_{i j}=\sum_{l=1}^{m} \mathbf{a}^{i j}\left(\xi_{l}^{2}-1\right)$, where $\mathbf{a}^{(i j)}:=\left(\left(X_{1}\right)_{i j}, \cdots,\left(X_{m}\right)_{i j}\right)^{T} \in \mathbb{R}^{m}$. Since $\xi_{l}$ is an i.i.d. copy of sub-Gaussian random variables, we know that there exist positive constants $M_{1}$ such that $\left\|\xi_{l}\right\|_{\psi_{1}} \leq M_{1}$ [42] Section 5.2.3]. Due to $\mathbb{E}\left(\xi_{l}\right)=0$, we know from [42] Proposition 5.10] that there exist positive constant $C_{4}$ such that for each $i, j \in\{1, \ldots, n\}$ and any given $t_{1}>0$,

$$
\mathbb{P}\left(\left|Y_{i j}\right| \geq t_{1}\right) \leq \exp \left(1-\frac{C_{4} t_{1}^{2}}{M_{1}^{2}\left\|\mathbf{a}^{(i j)}\right\|^{2}}\right)
$$

which implies that

$$
\begin{equation*}
\mathbb{P}\left(\|Y\|_{\infty} \geq t_{1}\right) \leq d_{\mathbb{S}_{h}^{n}} \exp \left(1-\frac{C_{4} t_{1}^{2}}{M_{1}^{2} \max \left\|\mathbf{a}^{(i j)}\right\|^{2}}\right) . \tag{46}
\end{equation*}
$$

Meanwhile, since $\xi_{l}$ is sub-Gaussian, we know that $\xi_{l}^{2}$ is an i.i.d. copy of sub-exponential random variables, which implies that there exists positive constant $M_{2}$ such that $\left\|\xi_{l}^{2}\right\|_{\psi_{1}} \leq M_{2}, l=1, \ldots, m$ (see e.g., [42], Section 5.2.4]). Moreover, since $\mathbb{E}\left(\xi_{l}^{2}\right)=1$, we know from [42, Proposition 5.16] that there exist positive constants $C_{5}$ such that for each $i, j \in\{1, \ldots, n\}$ and any given $t_{3}>0$,

$$
\mathbb{P}\left(\left|Z_{i j}\right| \geq t_{3}\right) \leq 2 \exp \left(-C_{5} \min \left\{\frac{t_{3}^{2}}{M_{2}^{2}\left\|\mathbf{a}^{(i j)}\right\|^{2}}, \frac{t_{3}}{M_{2}\left\|\mathbf{a}^{(i j)}\right\|_{\infty}}\right\}\right)
$$

which implies that

$$
\begin{equation*}
\mathbb{P}\left(\|Z\|_{\infty} \geq t_{3}\right) \leq 2 d_{\mathbb{S}_{h}^{n}} \exp \left(-C_{5} \min \left\{\frac{t_{3}^{2}}{M_{2}^{2} \max \left\|\mathbf{a}^{(i j)}\right\|^{2}}, \frac{t_{3}}{M_{2} \max \left\|\mathbf{a}^{(i j)}\right\|_{\infty}}\right\}\right) \tag{47}
\end{equation*}
$$

Moreover, for each $i, j \in\{1, \ldots, n\}$, it is clear that $\left\|\mathbf{a}^{(i j)}\right\|^{2} \leq m_{\max } / 4$ and $\left\|\mathbf{a}^{(i j)}\right\|_{\infty} \leq 1 / 2$, where $m_{\max }$ is the maximum number of repetition of $(i, j)$-th index in $\Omega$. Thus, it follows from 46) and 47) that for any given $t_{1}$, $t_{3}>0$,

$$
\begin{equation*}
\mathbb{P}\left(\|Y\|_{\infty} \geq t_{1}\right) \leq d_{\mathbb{S}_{h}^{n}} \exp \left(1-\frac{4 C_{4} t_{1}^{2}}{M_{1}^{2} m_{\max }}\right) \tag{48}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathbb{P}\left(\|Z\|_{\infty} \geq t_{3}\right) \leq 2 d_{\mathbb{S}_{h}^{n}} \exp \left(-C_{5} \min \left\{\frac{4 t_{3}^{2}}{M_{2}^{2} m_{\max }}, \frac{2 t_{3}}{M_{2}}\right\}\right) \tag{49}
\end{equation*}
$$

Therefore, it follows from Lemma 2 that with probability at least $1-1 /\left(2 n^{2}\right)$, there exists a constant $C_{6}>0$ such that $m_{\max } \leq C_{6} \log \left(2 n^{2}\right)$. Thus, by (48), we know that for any $t_{1}>0$,

$$
\begin{equation*}
\mathbb{P}\left(\left\|\mathscr{O}_{\Omega}^{*}(\xi)\right\|_{\infty} \geq t_{1}\right) \leq d_{\mathbb{S}_{h}^{n}} \exp \left(1-\frac{4 C_{4} t_{1}^{2}}{M_{1}^{2} C_{6} \log \left(2 n^{2}\right)}\right)+\frac{1}{2 n^{2}} \tag{50}
\end{equation*}
$$

On the other hand, since $\|Z\|_{\infty} \geq\left\|\mathscr{O}_{\Omega}^{*}(\xi \circ \xi)\right\|_{\infty}-\left\|\mathscr{O}_{\Omega}^{*}(\mathbf{1})\right\|_{\infty} \geq\left\|\mathscr{O}_{\Omega}^{*}(\xi \circ \xi)\right\|_{\infty}-m_{\max } / 2 \geq\left\|\mathscr{O}_{\Omega}^{*}(\xi \circ \xi)\right\|_{\infty}-$ $C_{6} \log \left(2 n^{2}\right) / 2$ if $m_{\max } \leq C_{6} \log \left(2 n^{2}\right)$, we know from Lemma 2 that for any $t_{2}>C_{6} \log \left(2 n^{2}\right) / 2$,

$$
\begin{align*}
\mathbb{P}\left(\left\|\mathscr{O}_{\Omega}^{*}(\xi \circ \xi)\right\|_{\infty} \geq t_{2}\right) & \leq \mathbb{P}\left(\|Z\|_{\infty} \geq t_{2}-\frac{C_{6} \log \left(2 n^{2}\right)}{2}\right)+\frac{1}{2 n^{2}} \\
& \leq 2 d_{\mathbb{S}_{h}^{n}} \exp \left(-C_{5} \min \left\{\frac{\left(2 t_{2}-C_{6} \log \left(2 n^{2}\right)\right)^{2}}{4 M_{2}^{2} C_{6} \log \left(2 n^{2}\right)}, \frac{2 t_{2}-C_{6} \log \left(2 n^{2}\right)}{M_{2}}\right\}\right)+\frac{1}{2 n^{2}} \tag{51}
\end{align*}
$$

Therefore, by setting $t_{1}:=M_{1} \sqrt{\frac{C_{6}}{2 C_{4}}} \log \left(2 n^{2}\right)$, we obtain from 50) that

$$
\mathbb{P}\left(\left\|\mathscr{O}_{\Omega}^{*}(\xi)\right\|_{\infty} \geq M_{1} \sqrt{\frac{C_{6}}{2 C_{4}}} \log \left(2 n^{2}\right)\right) \leq \frac{1}{2 n^{2}}+\frac{1}{2 n^{2}}=\frac{1}{n^{2}}
$$

Meanwhile, by setting $t_{2}:=\frac{4 M_{2}+1}{2} C_{6} \log \left(2 n^{2}\right)>C_{6} \log \left(2 n^{2}\right) / 2$, we conclude from 51) that

$$
\mathbb{P}\left(\left\|\mathscr{O}_{\Omega}^{*}(\xi \circ \xi)\right\|_{\infty} \geq \frac{4 M_{2}+1}{2} C_{6} \log \left(2 n^{2}\right)\right) \leq \frac{1}{2 n^{2}}+\frac{1}{2 n^{2}}=\frac{1}{n^{2}}
$$

Finally, it follows from (45) that there exists a constant $C_{3}>0$ such that

$$
\mathbb{P}\left(\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{\infty} \geq C_{3}\left(2 \omega \eta+\eta^{2}\right) \log \left(2 n^{2}\right)\right) \leq \frac{2}{n^{2}}
$$

which implies 44 holds with probability at least $1-2 / n^{2}$. This completes the proof.
Next, we shall present our statistical error bound results on the proposed convex model 12 . Proposition 4 and 5 suggest that the penalized parameters $\rho_{D}$ and $\rho_{S}$ based on 25 can take the following particular values:

$$
\begin{equation*}
\rho_{D}=O\left(\left(2 \omega \eta+\eta^{2}\right) \sqrt{\frac{\log (2 n)}{m n}}\right) \quad \text { and } \quad \rho_{S}=O\left(\left(2 \omega \eta+\eta^{2}\right) \frac{\log \left(2 n^{2}\right)}{m}\right) \tag{52}
\end{equation*}
$$

where $\omega=\left\|\mathscr{O}_{\Omega}(\bar{d}+\bar{s})\right\|_{\infty}$. Moreover, it follows from [11, (31)] and [46] Lemma 5.6] that if there exists $C_{1}>1$ such that $m>C_{1} n \log (n)$, then there exist positive constants $C_{4}$ and $C_{5}$ such that $\vartheta_{D}$ and $\vartheta_{S}$ defined by (35) satisfy

$$
\begin{equation*}
\vartheta_{D}=\mathbb{E}\left\|\frac{1}{m} \mathscr{O}_{\Omega}^{*}(\varepsilon)\right\|_{2} \leq C_{4} \sqrt{\frac{\log (2 n)}{m n}} \quad \text { and } \quad \vartheta_{S}=\mathbb{E}\left\|\frac{1}{m} \mathscr{O}_{\Omega}^{*}(\varepsilon)\right\|_{\infty} \leq C_{5} \frac{\log \left(2 n^{2}\right)}{m} \tag{53}
\end{equation*}
$$

Finally, by combining Proposition 3, 4 and 5, we obtain the following error bound, immediately. We omit the detail proof for the sake of brevity.

Theorem 1 Let $(\widehat{D}, \widehat{S})$ and $(\bar{D}, \bar{S})$ be an optimal solution of $(12)$ and the underground true EDM and outlier matrices, respectively. Assume the sample size $m$ satisfies $m>C_{1} n \log (2 n)$ for some constant $C_{1}>0$. For any given $\kappa_{D}>1$ and $\kappa_{S}>1$, suppose that the parameters $\rho_{D}>0$ and $\rho_{S}>0$ in the objective function (12) satisfy (52). Under Assumption 2. there exist some positive constants $C_{2}, C_{3}, C_{1}^{\prime}, C_{2}^{\prime}$ and $C_{3}^{\prime}$ such that either $\|\widehat{D}-\bar{D}\|^{2}+\|\widehat{S}-\bar{S}\|^{2} \leq \Gamma_{1}$ or

$$
\|\widehat{D}-\bar{D}\|^{2}+\|\widehat{S}-\bar{S}\|^{2} \leq \Gamma_{2}
$$

with probability at least $1-2 / n-2 / n^{2}$, where $\Gamma_{1}$ and $\Gamma_{2}$ are defined by

$$
\begin{equation*}
\Gamma_{1}:=C_{2} \mu_{1}\left(b_{D}+b_{S}\right)^{2} d_{\mathbb{S}_{h}^{n}} \sqrt{\frac{\log (2 n)}{m}} \tag{54}
\end{equation*}
$$

and

$$
\begin{align*}
\Gamma_{2}:= & C_{3} \mu_{1}^{2} d_{\mathbb{S}_{h}^{n}}\left\{C_{1}^{\prime} \eta^{2}(2 \omega+\eta)^{2}\left[\left(\kappa_{D} a_{D}+2 \sqrt{2}\right)^{2} \frac{r d_{\mathbb{S}_{h}^{n}} \log (2 n)}{n m}+\left(a_{S}+\frac{1}{\kappa_{S}}\right)^{2} \frac{k d_{\mathbb{S}_{h}^{n}} \log ^{2}\left(2 n^{2}\right)}{m^{2}}\right]\right. \\
& +C_{2}^{\prime}\left(b_{D}+b_{S}\right)^{2}\left(\frac{\kappa_{D}}{\kappa_{D}-1}\right)^{2}\left[\left(a_{D}+2 \sqrt{2}\right)^{2} \frac{r d_{h}^{n} \log (2 n)}{n m}+\left(\frac{\kappa_{S} a_{S}+1}{\kappa_{D}}\right)^{2} \frac{k d_{\mathbb{S}_{h}^{n}} \log ^{2}\left(2 n^{2}\right)}{m^{2}}\right] \\
& \left.+C_{3}^{\prime}\left(b_{D}+b_{S}\right)^{2}\left(\frac{\kappa_{S}}{\kappa_{S}-1}\right)^{2}\left[\left(\frac{\kappa_{D} a_{D}+2 \sqrt{2}}{\kappa_{S}}\right)^{2} \frac{r d_{\mathbb{S}_{h}^{n}} \log (2 n)}{n m}+\left(a_{S}+1\right)^{2} \frac{k d_{\mathbb{S}_{h}^{n}} \log ^{2}\left(2 n^{2}\right)}{m^{2}}\right]\right\}, \tag{55}
\end{align*}
$$

with $a_{D}$ and $a_{S}$ are given by (24).
We know from Theorem 1 that since the unknown true EDM $\bar{D}$ and outlier matrix $\bar{S}$ are bounded, in order to control the estimation error, we only need samples with the size $m$ of the order $\max \{r, k\}(n-1) \log (2 n) / 2$, since $d_{\mathbb{S}_{h}^{n}}=n(n-1) / 2$. Note that, it is reasonable to assume the embedding dimension $r=\operatorname{rank}(J \bar{D} J)$ and the outliers number $k$ are small. Therefore, the sample size $m$ is much smaller than $n(n-1) / 2$, the total number of the off-diagonal entries. However, we shall mention that one cannot obtain exact recovery from the bound obtained in Theorem 1 even without noise, i.e., $\eta=0$. Furthermore, as mentioned in [30], even for the outlier-free case (i.e., $\bar{S} \equiv 0$ ), this phenomenon is unavoidable due to lack of identifiability. For instance, consider the EDM $\bar{D}$ and the perturbed EDM $\widetilde{D}=\bar{D}+\varepsilon \mathbf{e}_{1} \mathbf{e}_{1}^{T}$. Thus, with high probability, $\mathscr{O}\left(D^{*}\right)=\mathscr{O}(\widetilde{D})$, which implies that it is impossible to distinguish two EDMs even if they are noiseless. If one is interested only in exact recovery in the noiseless setting, some addition assumptions such as the matrix incoherence condition (see e.g., [5, A0]) are necessary. In fact, recently, under matrix incoherence, random signs of outliers (i.e., the signs of the nonzero entries of $\bar{S}$ are i.i.d. symmetric Bernoulli random variables) and other assumptions, Chen et al. [8] obtained a near-optimal statistical guarantee of the convex nuclear norm plus $l_{1}$-norm penalized model for the (unconstrained) Robust PCA by building up the connection between the convex estimations and an auxiliary nonconvex optimization algorithm. For the Gaussian noise and squared matrices case, the estimation error bound achievable by their estimator [8] reads as

$$
\|\widehat{D}-\bar{D}\| \leq C \eta \sqrt{\frac{n}{m}}
$$

with high probability, where $C>0$ is a constant. Clearly, the resulting bound is stronger than ours for the case of the (unconstrained) Robust PCA. However, as we mentioned before, the results obtained in [8] have become inadequate since the model studied in their paper has no "hard-constraints", e.g., the noisy correlation matrix recovery (i.e., a positive semidefinite matrix whose diagonal elements are all ones) and the EDM estimation considered in this paper. Furthermore, neither matrix incoherence nor the random signs of outliers condition is assumed in this paper.

## 4 Recovery of the embedding dimensionality and outlier detection

In order to study the recovery guarantee of the dimensionality of embedding and outliers cardinality, we first introduce some useful notations and results on the proposed convex model (12) in Section 2 First, it is clear that the following generalized Slater condition for (12) always holds:

Definition 1 There exists $D^{0} \in \mathbb{S}^{n}$ and $S^{0} \in \mathbb{S}^{n}$ such that

$$
\operatorname{diag}\left(D^{0}\right)=0, \quad-D^{0} \in \operatorname{int}\left(\mathbb{K}_{+}^{n}\right) \quad \text { and } \quad S^{0}>0
$$

where $\mathbb{K}_{+}^{n}$ is the almost positive semidefinite matrix cone defined by $(2)$ and $\operatorname{int}\left(\mathbb{K}_{+}^{n}\right)$ is its interior.
Let $\left(\mathbb{K}_{+}^{n}\right)^{\circ} \subseteq \mathbb{S}^{n}$ be the polar cone of the almost positive semidefinite matrix cone $\mathbb{K}_{+}^{n}$, i.e.,

$$
\begin{equation*}
\left(\mathbb{K}_{+}^{n}\right)^{\circ}:=\left\{Z \in \mathbb{S}^{n} \mid\langle Z, Y\rangle \leq 0 \quad \forall Y \in \mathbb{K}_{+}^{n}\right\} \tag{56}
\end{equation*}
$$

We use $H \in \mathbb{S}^{n}$ to denote the Householder matrix, i.e.,

$$
\begin{equation*}
H:=I-\frac{2}{\mathbf{u}^{T} \mathbf{u}} \mathbf{u u}^{T} \quad \text { with } \quad \mathbf{u}:=(1, \ldots, 1, \sqrt{n}+1)^{T} \in \mathbb{R}^{n} . \tag{57}
\end{equation*}
$$

It is clear that the Householder matrix $H$ is symmetric and orthogonal (i.e., $H^{2}=I$ ). Also, the centering matrix $J$ defined by (1) satisfies

$$
J=H\left[\begin{array}{cc}
I_{n-1} & 0  \tag{58}\\
0 & 0
\end{array}\right] H .
$$

For any $X \in \mathbb{S}^{n}$, we rewrite the matrix $H X H$ as the following block form:

$$
H X H=\left[\begin{array}{cc}
\tilde{X}_{11} & \tilde{x}  \tag{59}\\
\tilde{x}^{T} & \tilde{x}_{0}
\end{array}\right] \quad \text { with } \quad \widetilde{X}_{11} \in \mathbb{S}^{n-1}, \quad \tilde{x} \in \mathbb{R}^{n-1} \quad \text { and } \quad \tilde{x}_{0} \in \mathbb{R} .
$$

Moreover, by (58) and simple calculations, we obtain the following basic identity:

$$
J X J=H\left[\begin{array}{rr}
\widetilde{X}_{11} & 0  \tag{60}\\
0 & 0
\end{array}\right] H
$$

where $\widetilde{X}_{11} \in \mathbb{S}^{n-1}$ is the first block defined by 59 for $H X H$.
By [17, Theorem 2.1], we have the following characterizations on $\mathbb{K}_{+}^{n}$ and its polar $\left(\mathbb{K}_{+}^{n}\right)^{\circ}$ :

$$
\mathbb{K}_{+}^{n}=\left\{\left.H\left[\begin{array}{cc}
Z & z  \tag{61}\\
z^{T} & z_{0}
\end{array}\right] H \in \mathbb{S}^{n} \right\rvert\, Z \in \mathbb{S}_{+}^{n-1}, z \in \mathbb{R}^{n-1}, z_{0} \in \mathbb{R}\right\}
$$

and

$$
\left(\mathbb{K}_{+}^{n}\right)^{\circ}=\left\{\left.H\left[\begin{array}{ll}
Z & 0  \tag{62}\\
0 & 0
\end{array}\right] H \in \mathbb{S}^{n} \right\rvert\, Z \in \mathbb{S}_{-}^{n-1}\right\}
$$

Thus, for any given integer $1 \leq r \leq n$, by (60) and 61), we know that $X \in \mathbb{K}_{+}^{n}$ and $\operatorname{rank}(J X J) \leq r$ if and only if $\widetilde{X}_{11} \in \mathbb{S}_{+}^{n-1}$ and $\operatorname{rank}\left(\widetilde{X}_{11}\right) \leq r$, and $\mathbb{K}_{+}^{n} \ni X \perp Y \in\left(\mathbb{K}_{+}^{n}\right)^{\circ}$ if and only if

$$
\mathbb{S}_{+}^{n-1} \ni \widetilde{X}_{11} \perp \widetilde{Y}_{11} \in \mathbb{S}_{-}^{n-1} \quad \text { and } \quad H Y H=\left[\begin{array}{cc}
\widetilde{Y}_{11} & 0  \tag{63}\\
0 & 0
\end{array}\right]
$$

where $\widetilde{X}_{11} \in \mathbb{S}^{n-1}$ and $\widetilde{Y}_{11} \in \mathbb{S}^{n-1}$ are the first blocks defined by (59) for $H X H$ and $H Y H$, respectively (see also [33] Lemma 2.1] for details).

Let $\widetilde{D} \in \mathbb{S}^{n}$ and $\widetilde{S} \in \mathbb{S}^{n}$ be the given initial estimators. Recall that $\widetilde{F} \in \mathbb{S}^{n}$ and $\widetilde{G} \in \mathbb{S}^{n}$ are the symmetric matrices defined in (14) and (17) with respect to $(\widetilde{D}, \widetilde{S})$. For the given $(\widetilde{D}, \widetilde{S})$, denote $\tilde{t}_{D} \in \mathbb{R}$ and $\tilde{t}_{S} \in \mathbb{R}$ by

$$
\tilde{t}_{D}=\left\{\begin{array}{ll}
\frac{\lambda_{r+1}\left(-\widetilde{D}_{11}\right)}{\lambda_{1}\left(-\widetilde{D}_{11}\right)} & \text { if } \widetilde{D} \neq 0,  \tag{64}\\
0 & \text { otherwise }
\end{array} \quad \text { and } \quad \tilde{t}_{S}= \begin{cases}\frac{\tilde{S}_{i^{\prime} j^{\prime}}}{\max _{k, l}\left\{\tilde{S}_{k}\right\}} & \text { if } \widetilde{S} \neq 0, \\
0 & \text { otherwise }\end{cases}\right.
$$

where $\left(i^{\prime}, j^{\prime}\right) \in\{1, \ldots, n\} \times\{1, \ldots, n\}$ be the index such that $\widetilde{S}_{i^{\prime} j^{\prime}}=\max \left\{\widetilde{S}_{i j} \mid(i, j) \notin \operatorname{supp}(\bar{S})\right\}$. Now, we are ready to present the results on the guarantee of recovery of the embedding dimensionality and outlier detection.

Theorem 2 Let $(\widehat{D}, \widehat{S})$ and $(\bar{D}, \bar{S})$ be an optimal solution of $\sqrt{12}$ and the underground true EDM and outlier matrices, respectively. Assume the sample size satisfies $m>C_{0} n \log (2 n)$ for some constant $C_{0}>0$. Suppose that the initial estimators $\widetilde{D}$ and $\widetilde{S}$ satisfy $\tilde{t}_{D} \in[0,1)$ and $\tilde{t}_{S} \in[0,1)$, where $\tilde{t}_{D}$ and $\tilde{t}_{S}$ are defined by (64). Let $\widetilde{F}$ and $\widetilde{G}$ be the symmetric matrices defined by (14) and (17) with respect to $(\widetilde{D}, \widetilde{S})$. Suppose that the parameters $\rho_{D}>0$ and $\rho_{S}>0$ in the objective function (12) defined by (52) satisfying $\rho_{D}>\frac{\varepsilon^{\tau}+\tilde{T}_{D}^{\tau}}{\varepsilon^{\tau}\left(1-\tilde{t}_{D}^{\tau}\right)} C\left(2 \omega \eta+\eta^{2}\right) \sqrt{\frac{\log (2 n)}{m n}}$ and $\rho_{S}>\frac{\varepsilon^{\tau}+\tilde{t}_{S}^{\tau}}{\varepsilon^{\tau}\left(1-\tau_{S}^{\tau}\right)} C\left(2 \omega \eta+\eta^{2}\right) \frac{\log \left(2 n^{2}\right)}{m}$ for some large constant $C>0$. Then, we have

$$
\operatorname{rank}(-J \widehat{D} J) \leq \operatorname{rank}(-J \bar{D} J) \quad \text { and } \quad \operatorname{supp}(\widehat{S}) \subseteq \operatorname{supp}(\bar{S})
$$

with probability at least $1-1 / n-3 / n^{2}$. Furthermore, in addition, if $\|\widehat{D}-\bar{D}\|<\lambda_{r}(-J \bar{D} J)$ and $\|\widehat{S}-\bar{S}\|<$ $\min \left\{\bar{S}_{i j} \mid(i, j) \in \operatorname{supp}(\bar{S})\right\}$, then with the same probability, we have

$$
\operatorname{rank}(-J \widehat{D} J)=\operatorname{rank}(-J \bar{D} J) \quad \text { and } \quad \operatorname{supp}(\widehat{S})=\operatorname{supp}(\bar{S})
$$

Proof Since (12) is convex and the generalized Slater condition (Definition 1) always holds, we know that there exist Lagrangian multipliers $(z, \Gamma, U) \in \mathbb{R}^{n} \times \mathbb{S}^{n} \times \mathbb{S}^{n}$ such that $(\widehat{D}, \widehat{S})$ satisfies the following Karush-Kuhn-Tucker (KKT) condition:

$$
\left\{\begin{array}{l}
-\frac{1}{m} \mathscr{O}_{\Omega}^{*}\left(\mathbf{y}-\mathscr{O}_{\Omega}(\widehat{D}+\widehat{S})\right)-\rho_{D} J(I-\widetilde{F}) J-\operatorname{Diag}(z)-\Gamma=0,  \tag{65}\\
-\frac{1}{m} \mathscr{O}_{\Omega}^{*}\left(\mathbf{y}-\mathscr{O}_{\Omega}(\widehat{D}+\widehat{S})\right)+\rho_{S}(E-\widetilde{G})+U=0, \\
\operatorname{diag}(\widehat{D})=0, \\
\mathbb{K}_{+}^{n} \ni-\widehat{D} \perp \Gamma \in\left(\mathbb{K}_{+}^{n}\right)^{\circ}, \quad 0 \leq \widehat{S}_{i j} \perp U_{i j} \leq 0, \quad i, j \in\{1, \ldots, n\} .
\end{array}\right.
$$

Consider the first equation of (65). By denoting $\Upsilon:=-\frac{1}{m} \mathscr{O}_{\Omega}^{*}\left(\mathbf{y}-\mathscr{O}_{\Omega}(\widehat{D}+\widehat{S})\right)$, we obtain that

$$
\begin{equation*}
H \operatorname{Diag}(z) H=H\left\ulcorner H-\rho_{D} H(J(I-\widetilde{F}) J) H-H \Gamma H\right. \tag{66}
\end{equation*}
$$

where $H$ is the Householder matrix defined by (57). Since $\Gamma \in\left(\mathbb{K}_{+}^{n}\right)^{\circ}$, we know from (60) and 62) that the last columns of the symmetric matrices $H(J(I-\widetilde{F}) J) H$ and $H \Gamma H$ are all zero. Moreover, for any $z \in \mathbb{R}^{n}$, we know that the last column of $H \operatorname{Diag}(z) H$ can be calculated as follows

$$
(H \operatorname{Diag}(z) H)_{(:, n)}=-\frac{1}{\sqrt{n}} H z \in \mathbb{R}^{n}
$$

Consequently, we know from (66) that the multiplier $z$ can be characterized by

$$
z=-\sqrt{n} H\left[\begin{array}{c}
\tilde{v} \\
\tilde{v}_{0}
\end{array}\right]
$$

where $\left(\tilde{v}, \tilde{v}_{0}\right)^{T} \in \mathbb{R}^{n}$ is the last column of $H \Upsilon H$ in the form $\left[59\right.$. Thus, since $\Upsilon \in \mathbb{S}_{h}^{n}$, we know from Lemma 1 that

$$
\|\operatorname{Diag}(z)\|_{2} \leq\|\operatorname{Diag}(z)\| \leq 2 \sqrt{n} \frac{1}{\sqrt{n}}\|r-J \Upsilon J\| \leq 2 \sqrt{2}\|r-J \Upsilon J\|_{2} \leq 4 \sqrt{2}\|r\|_{2}
$$

Meanwhile, by (10) and (7), we have

$$
\Upsilon=\frac{1}{m} \mathscr{O}_{\Omega}^{*}\left(\mathscr{O}_{\Omega}(\widehat{D}+\widehat{S})-\mathbf{y}\right)=\frac{1}{m} \mathscr{O}_{\Omega}^{*}\left(\mathscr{O}_{\Omega}\left(A_{D}+A_{S}\right)-\zeta\right),
$$

where $\left.\zeta=\mathscr{O}_{\Omega}( \rceil\right), A_{D}=\widehat{D}-\bar{D}$ and $A_{S}=\widehat{S}-\bar{S}$. Thus, we know that

$$
\|r\|_{2} \leq \frac{1}{m}\left\|\mathscr{O}_{\Omega}^{*} \mathscr{O}_{\Omega}\right\|_{2}\left\|A_{D}+A_{S}\right\|+\frac{1}{m}\left\|\mathscr{O}_{\Omega}^{*}(\zeta)\right\|_{2} .
$$

Under Assumption 11 we know that there exists a constant $C_{1}>0$ such that $\left\|A_{D}+A_{S}\right\| \leq C_{1}$. Moreover, by combining with Lemma2 Proposition 4 and $m>C_{0} n \log (2 n)$, we obtain that there exist positive constants $C_{2}>0$ and $C_{3}>0$ such that

$$
\|r\|_{2} \leq C_{1} C_{2} \frac{\log \left(2 n^{2}\right)}{m}+C_{3}\left(2 \omega \eta+\eta^{2}\right) \sqrt{\frac{\log (2 n)}{n m}} \leq C_{4}\left(2 \omega \eta+\eta^{2}\right) \sqrt{\frac{\log (2 n)}{n m}},
$$

with probability at least $1-1 /\left(2 n^{2}\right)-1 / n$, where $C_{4}:=C_{1} C_{2}+C_{3}>0$. Thus, by the assumption, we know that

$$
\begin{equation*}
\rho_{D}>(1+4 \sqrt{2}) \frac{\varepsilon^{\tau}+\tilde{t}_{D}^{\tau}}{\varepsilon^{\tau}\left(1-\tilde{t}_{D}^{\tau}\right)}\|r\|_{2} \tag{67}
\end{equation*}
$$

with probability at least $1-1 /\left(2 n^{2}\right)-1 / n$.
Denote $r=\operatorname{rank}(-J \bar{D} J)$. It is clear that $r \leq n-1$. By (14), we know that $\widetilde{F}$ satisfies

$$
J \widetilde{F} J=H\left[\begin{array}{rr}
\widetilde{\Lambda} & 0 \\
0 & 0
\end{array}\right] H \quad \text { with } \quad \tilde{\Lambda}=Q \operatorname{Diag}\left(f\left(\lambda\left(-\widetilde{D}_{11}\right)\right) Q^{T} \in \mathbb{S}^{n-1},\right.
$$

where $f: \mathbb{R}^{n-1} \rightarrow \mathbb{R}^{n-1}$ is the symmetric function defined by $(15)$, and $Q$ is a given $(n-1) \times(n-1)$ orthogonal matrix. It then follows from the well-known Weyl eigenvalue inequality [44] (see also [20, Theorem 4.3.7]) that

$$
\begin{align*}
\lambda_{r+1}\left(\widetilde{\Gamma}_{11}\right) & \leq \lambda_{1}(H(r-\operatorname{Diag}(z)) H)+\lambda_{r+1}\left(-\rho_{D}\left(I_{n-1}-\tilde{\Lambda}\right)\right) \\
& \leq\|r-\operatorname{Diag}(z)\|_{2}+\rho_{D} \lambda_{r+1}(\widetilde{\Lambda})-\rho_{D} \\
& \leq(1+4 \sqrt{2})\|r\|_{2}+\rho_{D}\left(\phi\left(\tilde{t}_{D}\right)-1\right) \tag{68}
\end{align*}
$$

where $\phi: \mathbb{R} \rightarrow \mathbb{R}$ is the scalar function given by (16), which implies that

$$
\lambda_{r+1}\left(\widetilde{\Gamma}_{11}\right) \leq(1+4 \sqrt{2})\|r\|_{2}+\rho_{D} \frac{\varepsilon^{\tau}\left(\tilde{t}_{D}^{\tau}-1\right)}{\tilde{t}_{D}^{\tau}+\varepsilon^{\tau}}
$$

Thus, we know from (67) that with probability at least $1-1 /\left(2 n^{2}\right)-1 / n, \lambda_{r+1}\left(\widetilde{\Gamma}_{11}\right)<0$. Since $\mathbb{K}_{+}^{n} \ni-\widehat{D} \perp \Gamma \in$ $\left(\mathbb{K}_{+}^{n}\right)^{\circ}$, we know from (63) that

$$
\begin{equation*}
\operatorname{rank}(-J \widehat{D} J) \leq \operatorname{rank}(-J \bar{D} J) \tag{69}
\end{equation*}
$$

with probability at least $1-1 /\left(2 n^{2}\right)-1 / n$.
Meanwhile, we know from the second equation of 65) that the multiplier $U \in \mathbb{S}^{n}$ is given by

$$
U=-\Upsilon-\rho_{S}(E-\widetilde{G})
$$

Again, under Assumption 1. we know that there exists a constant $C_{1}>0$ such that $\left\|A_{D}+A_{S}\right\| \leq C_{1}$. Moreover, by combining with Lemma 2 and Proposition 5 , we obtain that there exist positive constants $C_{2}$ and $C_{3}$ such that

$$
\|r\|_{\infty} \leq C_{1} C_{2} \frac{\log \left(2 n^{2}\right)}{m}+C_{3} \frac{\log \left(2 n^{2}\right)}{m} \leq C_{4} \frac{\log \left(2 n^{2}\right)}{m}
$$

with probability at least $1-5 /\left(2 n^{2}\right)$, where $C_{4}:=C_{1} C_{2}+C_{3}>0$. Therefore, by the assumption, we have

$$
\begin{equation*}
\rho_{S}>\frac{\varepsilon^{\tau}+\tilde{t}_{S}^{\tau}}{\varepsilon^{\tau}\left(1-\tilde{t}_{S}^{\tau}\right)}\|r\|_{\infty} \tag{70}
\end{equation*}
$$

with probability at least $1-5 /\left(2 n^{2}\right)$.
Let $(\hat{i}, \hat{j}) \in\{1, \ldots, n\} \times\{1, \ldots, n\}$ be the index such that $\widehat{S}_{\hat{i} \hat{j}}=\max \left\{\widehat{S}_{i j} \mid(i, j) \notin \operatorname{supp}(\bar{S})\right\}$. By (17), we know that

$$
U_{\hat{i} \hat{j}}=-\Upsilon_{\hat{i} \hat{j}}+\rho_{S}\left(\phi\left(\widetilde{S}_{\hat{i} \hat{j}} / \max _{k, l}\left\{\widetilde{S}_{k l}\right\}\right)-1\right)
$$

where $\phi: \mathbb{R} \rightarrow \mathbb{R}$ is the scalar function given by (16). It is clear from (64) that $\widetilde{S}_{\hat{i} \hat{j}} / \max _{k, l}\left\{\widetilde{S}_{k l}\right\} \leq \tilde{t}_{S}$. Therefore, since $\phi$ is non-decreasing, we have

$$
\begin{equation*}
U_{\hat{i} \hat{j}}=-\Upsilon_{\hat{i} \hat{j}}+\rho_{S}\left(\phi\left(\widetilde{S}_{\hat{i} \hat{j}} / \max _{k, l}\left\{\widetilde{S}_{k l}\right\}\right)-1\right) \leq-\Upsilon_{\hat{i} \hat{j}}+\rho_{S}\left(\phi\left(\tilde{t}_{S}\right)-1\right) \tag{71}
\end{equation*}
$$

Thus, we know from (71) and (16) that

$$
U_{\hat{i} \hat{j}} \leq\|r\|_{\infty}+\rho_{S}\left(\phi\left(\tilde{t}_{S}\right)-1\right)=\|r\|_{\infty}+\rho_{S} \frac{\varepsilon^{\tau}\left(\tilde{t}_{S}^{\tau}-1\right)}{\tilde{\tau}_{S}^{\tau}+\varepsilon^{\tau}}
$$

This, together with (70), yields $U_{\hat{i} \hat{j}}<0$ with probability at least $1-5 /\left(2 n^{2}\right)$. Moreover, since $0 \leq \widehat{S}_{i j} \perp U_{i j} \leq 0$ for any $i, j \in\{1, \ldots, n\}$, we know that with probability at least $1-5 /\left(2 n^{2}\right), \widehat{S}_{\hat{i} \hat{j}}=0$. By noting that for any $(i, j) \notin \operatorname{supp}(\bar{S}), 0 \leq \widehat{S}_{i j} \leq \widehat{S}_{\hat{i} \hat{j}}=0$, we conclude that

$$
\begin{equation*}
\operatorname{supp}(\widehat{S}) \subseteq \operatorname{supp}(\bar{S}) \tag{72}
\end{equation*}
$$

with probability at least $1-5 /\left(2 n^{2}\right)$.
By combining (69) and (72), we obtain that with probability at least $1-1 / n-3 /\left(n^{2}\right), \operatorname{rank}(-J \widehat{D} J) \leq \operatorname{rank}(-J \bar{D} J)$ and $\operatorname{supp}(\widehat{S}) \subseteq \operatorname{supp}(\bar{S})$. This completes the proof of the first part.

Next, we proceed with the proof of the second part. We know from the assumption $\|\widehat{D}-\bar{D}\|<\lambda_{r}(-J \bar{D} J)$ that

$$
\left|\lambda_{r}(-J \widehat{D} J)-\lambda_{r}(-J \bar{D} J)\right| \leq\|J \widehat{D} J-J \bar{D} J\| \leq\|\widehat{D}-\bar{D}\|<\lambda_{r}(-J \bar{D} J)
$$

which implies that $\lambda_{r}(-J \widehat{D} J)>0$. This yields

$$
\begin{equation*}
\operatorname{rank}(-J \widehat{D} J) \geq \operatorname{rank}(-J \bar{D} J) \tag{73}
\end{equation*}
$$

Meanwhile, let $(\bar{i}, \bar{j}) \in\{1, \ldots, n\} \times\{1, \ldots, n\}$ be the index such that $\bar{S}_{\bar{i} \bar{j}}=\min \left\{\bar{S}_{i j} \mid(i, j) \in \operatorname{supp}(\bar{S})\right\}$. Again, we know from the assumption $\|\widehat{S}-\bar{S}\|<\bar{S}_{\bar{i} \bar{j}}$ that

$$
\left|\widehat{S}_{i j}-\bar{S}_{i j}\right| \leq\|\widehat{S}-\bar{S}\|<\bar{S}_{\bar{i} \bar{j}} \quad \forall(i, j) \in \operatorname{supp}(\bar{S})
$$

This yields that for any $(i, j) \in \operatorname{supp}(\bar{S}), \widehat{S}_{i j}>\bar{S}_{i j}-\bar{S}_{\bar{i} \bar{j}} \geq 0$, which implies that

$$
\begin{equation*}
\operatorname{supp}(\widehat{S}) \supseteq \operatorname{supp}(\bar{S}) \tag{74}
\end{equation*}
$$

Therefore, by combining (73) and (74), we know from the first part of this theorem that with probability at least $1-1 / n-3 /\left(n^{2}\right), \operatorname{rank}(-J \widehat{D J})=\operatorname{rank}(-J \bar{D} J)$ and $\operatorname{supp}(\widehat{S})=\operatorname{supp}(\bar{S})$. The proof is completed.

Remark 1 In our implementations, we may choose the initial estimators $\widetilde{D}$ and $\widetilde{S}$ obtained by the nuclear norm $l_{1}$-minimization EDM problem (13) to generate $\widetilde{F}$ and $\widetilde{G}$ by $(14)$ and 17 , since the corresponding $\tilde{t}_{D}$ and $\tilde{t}_{S}$ satisfy $\tilde{t}_{D} \in[0,1)$ and $\tilde{t}_{S} \in[0,1)$ with high probability. Moreover, by combining Theorem 1 and Theorem 2 we know that if in addition $\lambda_{r}(-J \bar{D} J)>\max \left\{\Gamma_{1}^{1 / 2}, \Gamma_{2}^{1 / 2}\right\} \geq 0$ and

$$
\begin{equation*}
\min \left\{\bar{S}_{i j} \mid(i, j) \in \operatorname{supp}(\bar{S})\right\}>\max \left\{\Gamma_{1}^{1 / 2}, \Gamma_{2}^{1 / 2}\right\} \geq 0 \tag{75}
\end{equation*}
$$

where $\Gamma_{1}$ and $\Gamma_{2}$ are defined by (54) and (55), respectively, then

$$
\operatorname{rank}(-J \widehat{D} J)=\operatorname{rank}(-J \bar{D} J) \quad \text { and } \quad \operatorname{supp}(\widehat{S})=\operatorname{supp}(\bar{S})
$$

with probability at least $1-3 / n-5 / n^{2}$.

## 5 Numerical experiments

In this section, we shall demonstrate and verify the theoretical results obtained in Section 3 and 4 for the proposed matrix optimization model (12) by numerical experiments. In this paper, we directly employ the symmetric GaussSeidel decomposition based proximal alternating direction method of multipliers (sGS-ADMM) (cf. [38,25]) to solve the proposed matrix optimization model (12). The detail algorithm for solving (12] can be found in [10]. The numerical examples were tested on Matlab (2019b) under a Windows 10 64-bit Desktop (4 core, Intel Core i7-4790K @ $4.00 \mathrm{GHZ}, 16 \mathrm{~GB}$ RAM). We terminate sGS-ADMM if the KKT condition [10, (27)] are met, i.e.,

$$
\begin{equation*}
\max \left\{R_{p}, R_{d}, \text { rel_gap }\right\} \leq 10^{-4} \tag{76}
\end{equation*}
$$

where $R_{p}, R_{d}$ and rel_gap are the relative infeasibilities of the primal problem 12) and its dual problem, and the relative primal-dual gap, respectively, which are given by [10, (28)].

In order to demonstrate and verify the theoretical results, we only focus on the examples coming from a simulated network. For numerical performance results of the proposed model on real-world applications such as the no-line-sight mitigation in collaborative position localization, one may refer [10] for more details. Consider a randomly generated network in $\mathbb{R}^{r}$ with $r=2$, where $n$ points $\left\{p_{i}\right\} \in \mathbb{R}^{r}$ located randomly in the square area $[0,100] \times[0,100]$. We construct the observation operator $\mathscr{O}_{\Omega}$ defined in 9 by picking $\left\{X_{1}, \ldots, X_{m}\right\}$ uniformly at random from the standard basis matrices of the hellos space $\mathbb{S}_{h}^{n}$ with the sample size $m=O(r n \log (2 n))$. Meanwhile, we randomly add $k$ outliers which are modeled as the i.i.d. random variables to the true pairwise distances. The i.i.d. noise errors $\xi_{i j}$ in (5) follow a zero-mean Gaussian distribution with standard deviation and the noise magnitude control factor $\eta=0.5$. In all numerical experiments conducted in this paper, the parameters $\rho_{D}>0$ and $\rho_{S}>0$ in the convex model (12) are chosen exactly based on the rules suggested in (52), i.e., $\rho_{D}=O(\sqrt{\log (2 n) / m n})$ and $\rho_{S}=O\left(\log \left(2 n^{2}\right) / m\right)$. Meanwhile, the symmetric matrices $\widetilde{F}$ and $\widetilde{G}$ are defined by


Fig. 1: Performance comparison for different $\rho_{D}$ and $\rho_{S}$
(14) and (17) with respect to the initial estimators $\widetilde{D}$ and $\widetilde{S}$. In particular, we adopt the recommendation provided in [29, (25) and (26)] and [46, Chapter 5.3] with $\varepsilon \approx 0.05$ (within $0.01 \sim 0.1$ ), $\tau=2$ (within $1 \sim 3$ ) for the scalar function $\phi$ defined by (16). Also, the initial estimators $\widetilde{D}$ and $\widetilde{S}$ are generated from the nuclear norm and $l_{1}$ penalized least squares problem (i.e., the convex problem defined in (12) with $\widetilde{F}=0$ and $\widetilde{G}=0$ ). It seems that these particular settings for $\widetilde{F}, \widetilde{G}$ and initial estimators work quite well based on our numerical experiments.

Example 1. In this example, we use a simulation network with $n=1000$ random points to demonstrate the quality of the proposed estimators for different parameters $\rho_{D}$ and $\rho_{S}$. Here, $k=457$ outliers are modeled as the i.i.d. exponential random variables with the rate parameter $\lambda=100$. The sample size $m$ of the random observation operator $\mathscr{O}_{\Omega}$ equals to 15202 , which is in the order of $O(r n \log (2 n))$. The numerical performance for different parameters $\rho_{D}$ and $\rho_{S}$ are illustrated in Figure 1 The blue lines in both Figure 1a and 1b indicate the relative errors of both estimated EDM and outlier matrices, i.e., $\|\widehat{D}-\bar{D}\|+\|\widehat{S}-\bar{S}\|$, with respect to $\log \left(\rho_{D}\right)$ and $\rho_{S}$, respectively. It can be seen clearly when $\rho_{D}$ and $\rho_{S}$ are increasing, the relative errors are decreasing. The red line in Figure 1a stands for the estimated embedding dimension, i.e., $\operatorname{rank}(-J \widehat{D} J)$, which indicates that the estimated embedding dimension is always less or equal to the true $r=2$ when $\rho_{D}$ and $\rho_{S}$ are large enough. In fact, it actually equals $r=2$ for large $\rho_{D}$ and $\rho_{S}$. It is worth to note that in this example the $r$-th eigenvalue $\lambda_{r}(-J \bar{D} J)$ of $-J \bar{D} J$ is in the order of $O\left(10^{5}\right)$, which is much larger than $\max \left\{\Gamma_{1}^{1 / 2}, \Gamma_{2}^{1 / 2}\right\} \approx 10^{3}$, where $\Gamma_{1}$ and $\Gamma_{2}$ are the error bounds defined in (54) and 55), respectively. Meanwhile, the red line in Figure 1b represents the number of outlier detection errors obtained by the proposed convex model, which includes both false-negative errors denoted by FN (i.e., $\widehat{S}_{i j}=0$ but the true $\bar{S}_{i j}>0$ ) and false-positive errors denoted by FP (i.e., $\widehat{S}_{i j}>0$ but the true $\bar{S}_{i j}=0$ ). We know from Figure 1 b and Table 1 that the detected outlier number $\mathrm{nz}_{-} \mathrm{S}=456$, $\# \mathrm{FN}=1$ and $\# \mathrm{FP}=0$, when $\rho_{D}$ and $\rho_{S}$ are large enough. This implies that $\operatorname{supp}(\widehat{S}) \subseteq \operatorname{supp}(\bar{S})$, which is consistent with Theorem 2 Interestingly, in this example, we find that $\min \left\{\bar{S}_{i j} \mid(i, j) \in \operatorname{supp}(\bar{S})\right\} \approx 25.2695$ is much smaller that $\|\widehat{S}-\bar{S}\| \approx 10^{3}$. The detail numerical performance can be found in Table 1

Example 2. We use this example to illustrate the quality of the proposed estimators for problems with different dimensions. Ten simulation networks with $n=\{200,400, \ldots, 2000\}$ random points in the square area $[0,100] \times$ $[0,100]$ are generated in a similar manner as Example 1. The sample sizes of the corresponding observation operators and the numbers of outliers for different networks are reported in Table 2 The black lines in both Figure

| $\left(\rho_{D}, \rho_{S}\right)$ | $\left(R_{p}, R_{d}\right.$, rel_gap $)$ | $\left(r e l \_e r r, r \_d i m, n z \_\right.$S,mis_0) | cpu $(\mathrm{s})$ |
| :---: | :---: | :---: | :---: |
| $(0.0035,0.0194)$ | $(4.25 \mathrm{e}-06,7.18 \mathrm{e}-06,8.83 \mathrm{e}-05)$ | $(0.0064,31,2265,1808)$ | 360.84 |
| $(0.0040,0.0389)$ | $(6.35 \mathrm{e}-06,1.03 \mathrm{e}-05,8.09 \mathrm{e}-05)$ | $(0.0062,31,2323,2780)$ | 319.67 |
| $(0.0034,0.0865)$ | $(1.61 \mathrm{e}-06,2.35 \mathrm{e}-06,9.55 \mathrm{e}-05)$ | $(0.0060,30,797,342)$ | 321.04 |
| $(0.0052,0.0778)$ | $(1.28 \mathrm{e}-06,2.27 \mathrm{e}-06,9.71 \mathrm{e}-05)$ | $(0.0058,29,655,200)$ | 336.43 |
| $(0.0058,0.0972)$ | $(5.27 \mathrm{e}-07,5.07 \mathrm{e}-07,9.91 \mathrm{e}-05)$ | $(0.0056,29,468,13)$ | 430.75 |
| $(0.0288,0.1264)$ | $(3.72 \mathrm{e}-07,5.56 \mathrm{e}-07,9.96 \mathrm{e}-05)$ | $(0.0043,2,456,1)$ | 363.16 |
| $(0.0576,0.1361)$ | $(3.17 \mathrm{e}-07,4.35 \mathrm{e}-07,9.94 \mathrm{e}-05)$ | $(0.0042,2,456,1)$ | 304.29 |
| $(0.5762,0.1458)$ | $(9.41 \mathrm{e}-07,9.97 \mathrm{e}-07,6.69 \mathrm{e}-05)$ | $(0.0034,2,475,20)$ | 142.06 |
| $(2.8520,0.1555)$ | $(4.73 \mathrm{e}-06,8.29 \mathrm{e}-06,2.50 \mathrm{e}-06)$ | $(0.0022,2,529,74)$ | 328.50 |
| $(2.8808,0.1653)$ | $(5.91 \mathrm{e}-06,1.00 \mathrm{e}-05,3.18 \mathrm{e}-06)$ | $(0.0022,2,533,78)$ | 325.34 |
| $(5.7616,0.1750)$ | $(4.71 \mathrm{e}-07,3.85 \mathrm{e}-07,9.59 \mathrm{e}-05)$ | $(0.0017,2,456,1)$ | 795.77 |
| $(11.5233,0.1847)$ | $(1.12 \mathrm{e}-07,2.15 \mathrm{e}-08,9.73 \mathrm{e}-05)$ | $(0.0013,2,456,1)$ | 1357.65 |
| $(17.2849,0.1944)$ | $(2.52 \mathrm{e}-08,5.78 \mathrm{e}-09,9.93 \mathrm{e}-05)$ | $(0.0011,2,456,1)$ | 1890.97 |
| $(23.0465,0.2041)$ | $(1.57 \mathrm{e}-09,4.22 \mathrm{e}-09,9.98 \mathrm{e}-05)$ | $(0.0011,2,456,1)$ | 2401.03 |
| $(28.8082,0.2139)$ | $(6.43 \mathrm{e}-10,1.41 \mathrm{e}-09,1.00 \mathrm{e}-04)$ | $(0.0010,2,456,1)$ | 2783.24 |
| $(57.6164,0.2333)$ | $(5.60 \mathrm{e}-11,3.89 \mathrm{e}-11,9.98 \mathrm{e}-05)$ | $(0.0009,2,456,1)$ | 4079.01 |
| $(86.4245,0.2527)$ | $(1.76 \mathrm{e}-11,3.07 \mathrm{e}-11,9.98 \mathrm{e}-05)$ | $(0.0009,2,456,1)$ | 4581.41 |
| $(115.2327,0.2722)$ | $(2.94 \mathrm{e}-11,2.13 \mathrm{e}-11,9.98 \mathrm{e}-05)$ | $(0.0009,2,456,1)$ | 4592.07 |
| $(144.0409,0.2916)$ | $(1.73 \mathrm{e}-11,2.46 \mathrm{e}-11,9.99 \mathrm{e}-05)$ | $(0.0009,2,456,1)$ | 4633.12 |
| $(172.8491,0.311)$ | $(1.10 \mathrm{e}-11,1.47 \mathrm{e}-11,9.99 \mathrm{e}-05)$ | $(0.0009,2,456,1)$ | 4809.25 |

Table 1: Numerical performance of the network with $(n, m, k)=(1000,15202,457)$ for different $\rho_{D}$ and $\rho_{S}: R_{p}$, $R_{d}$ and rel_gap stand for the relative primal feasibility, dual feasibility and relative duality gap obtained by sGS-ADMM, respectively; we use rel_err, r_dim, nz_S, mis_0 to denote the relative error, recovery embedding dimension, detected outlier number and outlier detection error; cpu ( $s$ ) is the total computational time (in seconds) of sGS-ADMM.

2 a and 2 b represent the theoretical (relative) upper bounds in Theorem 1 , i.e., $\max \left(\Gamma_{1}, \Gamma_{2}\right) /(1+\|\bar{D}\|+\|\bar{S}\|)$ and $\Gamma_{1}$ and $\Gamma_{2}$ are defined in (54) and 555 , respectively. The blue lines in both Figure 2 a and 2 b are the square sum of relative errors with respect to $\widehat{D}$ and $\widehat{S}$, i.e., $\|\widehat{D}-\bar{D}\|^{2}+\|\widehat{S}-\bar{S}\|^{2} /(1+\|\bar{D}\|+\|\bar{S}\|)$. It can be seen clearly from Figure 2 a and 2 b that the square sum of (relative) errors with respect to $\widehat{D}$ and $\widehat{S}$ is smaller than theoretical (relative) upper bounds defined by (54) and (55). It can be seen from Figure 2athat the proposed convex model (12) provides the estimators $\widehat{D}$ with the true EDM dimension, i.e., $\operatorname{rank}(-J \widehat{D} J)=r=2$ in all ten networks. Similar with Example 1, we note that for these ten cases, the $r$-th eigenvalues $\lambda_{r}(-J \bar{D} J)$ of $-J \bar{D} J$ are in the order of $O\left(10^{5}\right)$, and the upper bounds defined in (54) and (55) satisfy $\max \left\{\Gamma_{1}^{1 / 2}, \Gamma_{2}^{1 / 2}\right\} \approx 10^{3}$. One the other hand, the red line in Figure 2bindicates the number of outlier detection errors obtained by 12 for different networks. We know from Figure 2 b and Table 2 that for all networks with different dimension scales, the detected outlier numbers $n z_{-} \leq \leq k$ and nz_S + mis_ $0=k$, which implies that for each case, the outlier detection errors if exist are the false-negative errors (FN) and $\operatorname{supp}(\widehat{S}) \subseteq \operatorname{supp}(\bar{S})$. Similarly with Example 1, this result is consistent with the outlier detection guarantee results proposed in Theorem 2 Also, it is worth noting that $10 \approx \min \left\{\bar{S}_{i j} \mid(i, j) \in \operatorname{supp}(\bar{S})\right\} \ll \| \widehat{S}-$ $\bar{S} \| \approx 10^{3}$ for these cases. The numerical details are reported in Table 2

Example 3. Finally, we conduct an experiment to verify the proposed sufficient condition for the outlier detection in Theorem 2 Consider a simulation network with $n=200$ points which are randomly located in the


Fig. 2: Performance comparison for networks with different dimension scales

| $(n, m, k)$ | $\left(\rho_{D}, \rho_{S}\right)$ | $\left(R_{p}, R_{d}, g a p\right)$ | $($ rel_err,r_dim,nz_S,mis_0) | cpu $(\mathrm{s})$ |
| :---: | :---: | :---: | :---: | :---: |
| $(200,2397,72)$ | $(393.7605,1.5738)$ | $(8.44 \mathrm{e}-11,7.14 \mathrm{e}-11,9.95 \mathrm{e}-05)$ | $(0.0018,2,72,0)$ | 17.30 |
| $(400,5348,161)$ | $(157.6384,0.6341)$ | $(2.67 \mathrm{e}-11,7.02 \mathrm{e}-11,9.97 \mathrm{e}-05)$ | $(0.0015,2,1610)$ | 193.20 |
| $(600,8509,256)$ | $(157.3630,0.6350)$ | $(3.04 \mathrm{e}-11,4.36 \mathrm{e}-11,1.00 \mathrm{e}-04)$ | $(0.0010,2,255,1)$ | 688.16 |
| $(800,11805,355)$ | $(111.6860,0.4516)$ | $(4.87 \mathrm{e}-11,3.19 \mathrm{e}-11,9.99 \mathrm{e}-05)$ | $(0.0010,2,354,1)$ | 1974.85 |
| $(1000,15202,457)$ | $(125.0299,0.5063)$ | $(1.92 \mathrm{e}-11,2.01 \mathrm{e}-11,1.00 \mathrm{e}-04)$ | $(0.0008,2,457,0)$ | 4583.55 |
| $(1200,18680,561)$ | $(69.6896,0.2825)$ | $(2.70 \mathrm{e}-11,2.87 \mathrm{e}-11,9.99 \mathrm{e}-05)$ | $(0.0010,2,561,0)$ | 7161.07 |
| $(1400,2225,667)$ | $(88.5303,0.3592)$ | $(1.96 \mathrm{e}-11,1.71 \mathrm{e}-11,9.99 \mathrm{e}-05)$ | $(0.0008,2,665,2)$ | 12641.73 |
| $(1600,25827,775)$ | $(78.9857,0.3207)$ | $(1.49 \mathrm{e}-11,1.57 \mathrm{e}-11,1.00 \mathrm{e}-04)$ | $(0.0008,2,774,1)$ | 21380.81 |
| $(1800,29480,885)$ | $(52.1488,0.2119)$ | $(1.69 \mathrm{e}-11,2.11 \mathrm{e}-11,1.00 \mathrm{e}-04)$ | $(0.0008,2,883,2)$ | 33169.30 |
| $(2000,33177,996)$ | $(44.9024,0.1825)$ | $(2.24 \mathrm{e}-11,1.93 \mathrm{e}-11,9.99 \mathrm{e}-05)$ | $(0.0008,2,993,3)$ | 48297.41 |

Table 2: Numerical performance for networks with different dimension scales: $n, m, k$ indicate the number of points in network, sample size of the observation operator and number of outliers; $R_{p}, R_{d}$ and rel_gap stand for the relative primal feasibility, dual feasibility and relative duality gap obtained by sGS-ADMM, respectively; we use rel_err, r_dim, nz_S, mis_0 to denote the relative error, recovery embedding dimension, detected outlier number and outlier detection error; $\mathrm{cpu}(\mathrm{s})$ is the total computational time (in seconds) of sGS-ADMM.
area $[0,100] \times[0,100]$. The observation operator $\mathscr{O}_{\Omega}$ and i.i.d. noise errors $\xi_{i j}$ are generated in the same manner as those in Example $1 \& 2$. Moreover, we randomly add $k=24$ outliers errors, which are the i.i.d. uniform random variables with different magnitudes such that $\bar{S}$ satisfies one of the following conditions, respectively: (a) $\min \left\{\bar{S}_{i j} \mid(i, j) \in \operatorname{supp}(\bar{S})\right\} \approx 1 \times 10$ denoted by the "small magnitude" of $\bar{S},(\mathrm{~b}) \min \left\{\bar{S}_{i j} \mid(i, j) \in \operatorname{supp}(\bar{S})\right\} \approx 10^{2}$ denoted by the "middle magnitude" of $\bar{S}$, and (c) $\min \left\{\bar{S}_{i j} \mid(i, j) \in \operatorname{supp}(\bar{S})\right\} \approx 5 \times 10^{3}$ denoted by the "large magnitude" of $\bar{S}$. Note that in this example, we have $\max \left\{\Gamma_{1}^{1 / 2}, \Gamma_{2}^{1 / 2}\right\} \approx 10^{3}$. By (75), we know that the condition $\|\widehat{S}-\bar{S}\|<\min \left\{\bar{S}_{i j} \mid(i, j) \in \operatorname{supp}(\bar{S})\right\}$ is satisfied in the "large magnitude" case. First, the estimators $\widehat{D}$ obtained by the convex model (12) in all cases satisfy $\operatorname{rank}(-J \widehat{D} J)=r=2$, since the $r$-th eigenvalues $\lambda_{r}(-J \bar{D} J)$ of $-J \bar{D} J$ are in the order of $O\left(10^{5}\right)$. For comparison, we report the relative errors and the outlier detection errors


Fig. 3: Relative error and outlier detection error for different magnitudes of $\bar{S}$
after 100 Monte Carlo simulation runs in Figure 3 Note that for all cases, the outlier detection errors if exist are the false-negative errors $(\mathrm{FN})$ and $\operatorname{supp}(\widehat{S}) \subseteq \operatorname{supp}(\bar{S})$, since the detected outlier numbers $\mathrm{nz}_{-} \mathrm{S} \leq k=$ and nz _S + mis_ $0=k$. For the small and middle magnitude of $\bar{S}$ cases (Figure 3a and 3b, only a few estimators $\widehat{S}$ satisfy $\operatorname{supp}(\widehat{S})=\operatorname{supp}(\bar{S})$ exactly ( 13 out of 100 MC simulation runs for the small magnitude case; 77 out of 100 MC simulation runs for the middle magnitude case). However, it is worth noting that for the large magnitude of $\bar{S}$ case (Figure 3c), with probability 1 , the estimator $\widehat{S}$ obtained by (12) satisfies supp $(\widehat{S})=\operatorname{supp}(\bar{S})$ exactly $(100$ out of 100 MC simulation runs). Also the empirical cumulative distribution function ${ }^{11}$ (CDF) of different cases are reported in Figure 3d

[^1]
## 6 Conclusions

Euclidean embedding from noisy observations containing outliers is an important and challenging problem in statistics and machine learning. Many existing methods would struggle with outliers due to a lack of detection ability, while the matrix optimization based embedding model introduced in [10] usually can produce reliable embeddings and identify the outliers jointly. This paper aimed to explain this mysterious situation by studying the estimation error bounds and outliers detection ability of the proposed model. In particular, we show that the estimators obtained by the proposed method satisfy a non-asymptotic risk bound, implying that the model provides a high accuracy estimator with high probability when the order of the sample size is roughly the degree of freedom up to a logarithmic factor. Moreover, we show that under some mild conditions, the proposed model also can identify the outliers without any prior information with high probability. As we mentioned in the Section 3. Chen, et al. [8], derived a near-optimal statistical guarantee of the convex nuclear norm plus $l_{1}$-norm penalized model for the classical Robust PCA by building up the connection between the convex estimations and an auxiliary nonconvex optimization algorithm. It seems their approach would lead to some interesting error bound results for our EDM embedding model. However, it seems difficult to extend their results directly to the convex models involving "hard-constraints", e.g., the noisy correlation matrix recovery and the EDM estimation studied in this paper. We plan to investigate those issues in the future.

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[^1]:    ${ }^{1}$ Let $x_{1}, \ldots, x_{n}$ be independent, identically distributed real random variables. The corresponding empirical distribution function $F_{n}(t)$ is defined as $F_{n}(t)=\frac{1}{n} \sum_{i=1}^{n} \delta_{x_{i} \leq t}$, where $\delta_{x_{i} \leq t}$ is the indicator of event $x_{i} \leq t$.

