M-estimation in Low-rank Matrix Factorization: a General Framework

Peng Liu1,* Wei Tu2,* Jingyu Zhao3 Yi Liu2 Linglong Kong2,† Guodong Li3 Bei Jiang2 Hengshuai Yao4
Guangjian Tian5
1School of Mathematics, Statistics and Actuarial Science, University of Kent
2Department of Mathematical and Statistical Sciences, University of Alberta
3Department of Statistics and Actuarial Science, University of Hong Kong
4Huawei Hi-Silicon, Canada
5Huawei Noah’s Ark Lab, Hong Kong, China
Emails: p.liu@kent.ac.uk.
{wei.tu, yliu16, lkong, bei1}@ualberta. {gladys17, gdli}@hku.hk. {hengshuai.yao, tian.guangjian}@huawei.com

Abstract—Many problems in science and engineering can be reduced to the recovery of an unknown large matrix from a small number of random linear measurements. Matrix factorization arguably is the most popular approach for low-rank matrix recovery. Many methods have been proposed using different loss functions, such as the most widely used $L_2$ loss, more robust choices $L_1$ and Huber loss, and quantile and expectile loss for skewed data. All of them can be unified into the framework of M-estimation. In this paper, we present a general framework of low-rank matrix factorization based on M-estimation in statistics. The framework mainly involves two steps: we first apply Nesterov’s smoothing technique to obtain an optimal smooth approximation for non-smooth loss functions, such as $L_1$ and quantile loss; secondly, we exploit an alternative updating scheme along with Nesterov’s momentum method at each step to minimize the smoothed loss function. Strong theoretical convergence guarantee has been developed for the general framework, and extensive numerical experiments have been conducted to illustrate the performance of the proposed algorithm.

Index Terms—matrix recovery, M-estimation, matrix factorization, robustness, statistical foundation

I. INTRODUCTION

Motivation. In matrix recovery from linear measurements, we are interested in recovering an unknown matrix $X \in \mathbb{R}^{m \times n}$ from $p < mn$ linear measurements $b_i = \text{Tr}(A_i^T X)$, where each $A_i \in \mathbb{R}^{m \times n}$ is a measurement matrix, $i = 1, \ldots, p$. Usually it is expensive or even impossible to fully sample the entire matrix $X$, and we are left with a highly incomplete set of observations. In general it is not always possible to recover $X$ under such settings. However, if we impose a low-rank structure on $X$, it is possible to exploit this structure and efficiently estimate $X$. The problem of matrix recovery arises in a wide range of applications, such as collabora filter [1], image recovery [2], structure from motion, photometric stereo [3], system identification [4] and computer network tomography [5].

*These authors contributed equally to this work.
†Correspondence Author. Part of the work is done during a sabbatical at Huawei.

Matrix factorization arguably is the most popular and intuitive approach for low-rank matrix recovery. The basic idea is to decompose the low-rank matrix $X \in \mathbb{R}^{m \times n}$ into the product of two matrices

$$X = U^T V,$$

where $U \in \mathbb{R}^{r \times m}$ and $V \in \mathbb{R}^{r \times n}$. In many practical problems, the rank of a matrix is known or can be estimated in advance; see, for example, the rigid and nonrigid structures from motion as well as image recovery. The matrix $U$ and $V$ can also be interpreted as latent factors that drive the unknown matrix $X$.

Matrix factorization is usually based on $L_2$ loss, i.e. square loss, which is optimal for Gaussian errors. However, its performance may be severely deteriorated when the data is contaminated by outliers. For example, in a collaborative filtering system, some popular items have a lot of ratings regardless of whether they are useful, while others have fewer ones. There may even exist shilling attacks, i.e. a user may consistently give positive feedback to their products or negative feedback to their competitors regardless of the items themselves [6]. Recently there are a few attempts to address this problem; see, for example, [2], [7]. However, to the best of our knowledge, they focus on either $L_1$ or quantile loss, and are only useful in limited scenarios. For low-rank matrix recovery under M-estimation, He et al. [8] studied the use of a few smooth loss functions such as Huber and Welsch in this setting. In this paper, we propose a more general framework that is applicable to any M-estimation loss function, smooth or non-smooth, and provide theoretical convergence guarantee on the proposed state-of-the-art algorithm.

This paper introduces a general framework of M-estimation [9]–[11] on matrix factorization. Specifically, we consider loss functions related to M-estimation for matrix factorization. M-estimation is defined in a way similar to the well-known terminology “Maximum Likelihood Estimate (MLE)” in statistics, and has many good properties that are similar to those of MLE. Meanwhile, it still retains intuitive interpretation. The proposed class of loss functions includes well-known $L_1$ and
The matrix $L_2$ loss as special cases. In practice, we choose a suitable M-estimation procedure according to our knowledge of the data and the specific nature of the problem. For example, Huber loss enjoys the property of smoothness as $L_2$ loss and robustness as $L_1$ loss.

The loss functions of some M-estimation procedures are smooth, such as the $L_2$ loss, while some others are non-smooth such as $L_1$ and quantile loss. Note that the resulting objective functions all have a bilinear structure due to the decomposition at Eq. (1). For non-smooth cases, we first consider Nesterov’s smoothing method to obtain an optimal smooth approximation [12], and the bilinear structure is preserved. The alternating minimization method is hence used to search for the solutions.

At each step, we employ Nesterov’s momentum method to accelerate the convergence, and it turns out to find the global optima easily. Figure 1 gives the flowchart of the proposed algorithm. Theoretical convergence analysis is conducted for both smooth and non-smooth loss functions.

**Contributions.** We summarize our contributions below.

1. We propose to do matrix factorization based on the loss functions of M-estimation procedures. The proposed framework is very general and applicable to any M-estimate loss function, which gives us flexibility in selecting a suitable loss for specific problems.

2. We propose to use Nesterov’s smoothing technique to obtain an optimal smooth approximation when the loss function is non-smooth.

3. We consider the Nesterov’s momentum method, rather than gradient descent methods, to perform the optimization at each step of the alternating minimization, which greatly accelerates the convergence.

4. We provide theoretical convergence guarantees for the proposed algorithm.

5. We illustrate the usefulness of our method by conducting extensive numerical experiments on both synthetic data and real data.

II. METHODOLOGY FRAMEWORK

Let $X^* \in \mathbb{R}^{m \times n}$ be the target low-rank matrix, and $A_i \in \mathbb{R}^{m \times n}$ with $1 \leq i \leq p$ be given measurement matrices. Here $A_i$’s can be the same or different with each other. We assume the observed signals $b = (b_1, ..., b_p)^T$ have the structure

$$b_i = \langle A_i, X^* \rangle + \epsilon_i, \quad i = 1, ..., p, \quad (2)$$

where $\langle A_i, X \rangle := \text{Tr}(A_i^T X)$ and $\epsilon_i$ is the error term. Suppose that the rank of matrix $X^*$ is no more than $r$ with $r \ll \min(m, n, p)$. We then have the decomposition $X^* = U^{r \times r} V^r$, and the matrix can be recovered by solving a non-convex optimization problem

$$\min_{U \in \mathbb{R}^{r \times m}, V \in \mathbb{R}^{r \times n}} \frac{1}{p} \sum_{i=1}^{p} \mathcal{L}(b_i - \langle A_i, U^T V \rangle), \quad (3)$$

where $\mathcal{L}(\cdot)$ is the loss function used an M-estimation procedure. For example, $\mathcal{L}(y) = y^2$ for the $L_2$ loss and $\mathcal{L}(y) = |y|$ for the $L_1$ loss. Here $\mathcal{L}(\cdot)$ usually is convex. Let $\mathcal{A} : \mathbb{R}^{m \times n} \rightarrow \mathbb{R}^p$ be an affine transformation with the $i$th entry of $\mathcal{A}(X)$ being $\langle A_i, X \rangle$, and $\mathcal{M}(x) = p^{-1} \sum_{i=1}^{p} \mathcal{L}(x_i)$ for a vector $x = (x_1, ..., x_p)^T$. We can then rewrite (3) into a more compact form

$$\min_{U \in \mathbb{R}^{r \times m}, V \in \mathbb{R}^{r \times n}} \mathcal{M}(b - \mathcal{A}(U^T V)). \quad (4)$$

A. The case that $\mathcal{M}$ is not smooth

The loss function $\mathcal{L}(\cdot)$, and hence $\mathcal{M}(\cdot)$, may be non-smooth in some M-estimation procedures, such as the $L_1$ loss. For this non-smooth case, we first employ Nesterov’s smoothing method to obtain an optimal smooth approximation; see [12] (pp. 129-132).

Specifically, we first assume that the objective function $\mathcal{M}$ has the following structure

$$\mathcal{M}(b - \mathcal{A}(U^T V)) = \hat{\mathcal{M}}(b - \mathcal{A}(U^T V))$$

$$+ \max_u \left\{ \langle B(b - \mathcal{A}(U^T V)), u \rangle - \phi(u) \right\}, \quad (5)$$

where $\hat{\mathcal{M}}(\cdot)$ is continuous and convex; see Eq. (2.2) in [12]. Then the objective function $\mathcal{M}$ can be approximated by

$$\mathcal{M}_\pi(b - \mathcal{A}(U^T V)) = \hat{\mathcal{M}}(b - \mathcal{A}(U^T V))$$

$$\max_u \left\{ \langle B(b - \mathcal{A}(U^T V)), u \rangle - \phi(u) - \pi d_2(u) \right\}, \quad (6)$$

where $\pi$ is a positive smoothness parameter. Note that $\mathcal{M}_\pi(\cdot)$ is smooth and can be optimized by many available gradient-based methods.
where the initialization step, we do not need a very small tolerance \( \epsilon \). Similarly define the regularized objective function, denoted by \( \mathcal{M} \), and all algorithms designed below are for this optimization below [13]. Specifically, rather than Eq. (4), we attempt to do the unstable solutions, and this paper solves this problem by keeping one of \( U \) and \( V \) fixed, optimize over the other, and then switch in the next iteration until the algorithm converges. Moreover, direct optimization of Eq. (4) may lead to solution, and this paper solves this problem by keeping one of \( U \) and \( V \) fixed, optimize over the other, and then switch in the next iteration until the algorithm converges.

### III. Initialization

When the starting values of \( U \) and \( V \) are orthogonal (or almost orthogonal) to the true space, an alternating minimization algorithm may never converge to true values, and hence an initialization procedure is needed to avoid this situation. Here we adopt the singular value projection (SVP), which originates from [14] and was later used by [13] to provide the initial values of \( U \) and \( V \) for the designed algorithm in the next section. The main difference between our method and the original one is summarized into Line 2 of Algorithm 1, where we use a loss function related to M-estimation,

\[
\nabla_X \mathcal{M}(b - \mathcal{A}(X^t)) = - \frac{1}{p} \sum_{i=1}^{p} \hat{L}(b - \mathcal{A}(X^t)) A_i,
\]

where \( \hat{L}(\cdot) \) is the first derivative of \( L(\cdot) \).

#### Algorithm 1: Initialization by SVP algorithm

**Input:** \( A, b, \) tolerance \( \epsilon_1, \) step size \( \xi_t \) with \( t = 0, 1, \cdots, \) and \( X^0 = 0_{m \times n} \)

**Output:** \( X^{t+1} \)

1. **Repeat**
   2. \( Y^{t+1} \leftarrow X^t - \xi_t \nabla_X \mathcal{M}(b - \mathcal{A}(X^t)) \)
   3. Compute top \( r \) singular vectors of \( Y^{t+1} \): \( U_r, \Sigma_r, V_r \)
   4. \( X^{t+1} \leftarrow U_r \Sigma_r V_r \)
   5. \( t \leftarrow t + 1 \)
   6. Until \( \|X^{t+1} - X^t\|_F \leq \epsilon_1 \)

It is noteworthy to point out that Algorithm 1 can be directly used to recover the matrix \( X^* \) if it is iterated for sufficient times. However, the singular value calculation here is time-consuming when the dimension of matrix \( X^* \) is large. In this initialization step, we do not need a very small tolerance \( \epsilon_1 \), i.e., a rough output is sufficient. Our simulation experiments show that, after several iterations of the SVP algorithm, the resulting values are close to the true ones, while it is not the case for random initialization.

Algorithm 1 can be rewritten into a compact form

\[
X^{t+1} \leftarrow \mathcal{P}_r \left( X^t - \xi_t \nabla_X \mathcal{M}(b - \mathcal{A}(X^t)) \right),
\]

where \( \mathcal{P}_r \) denotes the operation of projection onto the space of rank-\( r \) matrices. Moreover, the original objective function \( \mathcal{M} \), rather than the regularized one at Eq. (7), is used in Algorithm 1, and for the non-smooth case we will use \( \mathcal{M}_r \) at Eq. (6).

#### IV. Algorithm

There are two layers of iterations in our algorithm: the outer layer is the alternating minimization procedure; and the inner layer employs Nesterov’s momentum algorithm to obtain updated values of \( U^{t+1} \) or \( V^{t+1} \).

#### Algorithm 2: Nesterov’s accelerate gradient (NAG) method

**Input:** \( U^t, V^t, \) momentum parameter \( \gamma \), learning rate \( \eta \), and tolerance \( \epsilon_2 \)

**Output:** \( U^{t+1} \)

1. **Repeat**
   2. \( \nu_i^{t} = \gamma \nu_i^{t-1} + \eta \nabla_U \mathcal{M}(U_i^{t-1} - \gamma \nu_i^{t-1}, V^t) \)
   3. \( U_i^{t} = U_i^{t-1} + \nu_i^{t} \)
   4. Until \( \|U_i^{t} - U_i^{t-1}\|_F \leq \epsilon_2 \)

We first introduce the inner layer, where the Nesterov’s momentum method is used to update the values of \( U^t \) and \( V^t \) to those of \( U^{t+1} \) and \( V^{t+1} \). Algorithm 2 gives the details of updating the value of \( U^t \), and we denote it by NAGU for simplicity. Here the gradient \( \nabla_U \mathcal{M}(U, V) \) is defined as the gradient with respect to \( U \), and similarly we can define \( \nabla_V \mathcal{M}(U, V) \). Moreover, \( \nu_i^{t} \) stands for the momentum term, \( \gamma \) is the momentum parameter, and \( \eta \) is the learning rate. The value of \( \gamma \) is usually chosen to be around 0.9; see [15] and [16]. Similarly, we can give the detailed algorithm for updating the value of \( V^t \), and it can be denoted by NAGV.

The alternating minimization method is employed for the outer layer of iterations; see Algorithm 3 for details. The final solutions can be denoted by \( \hat{U} \) and \( \hat{V} \), and we then can use \( \hat{U}^\top \hat{V} \) to approximate the low-rank matrix \( X^* \).

#### Algorithm 3: Alternating Minimization

**Input:** \( U^0, V^0 \)

**Output:** \( \hat{U}, \hat{V} \)

1. **Repeat**
   2. Update \( U^t \) with \( U_i^{t+1} = \text{NAGU}(U_i^t, V^t) \)
   3. Update \( V^t \) with \( V_i^{t+1} = \text{NAGV}(U_i^{t+1}, V^t) \)
   4. Until converge
V. CONVERGENCE RESULTS

Suppose that $U$ and $V$ are a pair of solutions, i.e. $X = U^TV$. It then holds that, for an orthonormal matrix $R$ satisfying $R^TR = I_r$, $U^R = RU$ and $V^R = RV$ are another pair of solutions. To evaluate the performance of the proposed algorithm, we first define a distance between two matrices

$$\text{dist}(U, U^R) = \min_{R \in \mathbb{R}^{r \times r}, R^TR = I_r} \|U - RU^R\|_F,$$

where $U, U^R \in \mathbb{R}^{r \times m}$ with $m \geq r$; see [13].

Definition 1. (Restricted Isometry Property (RIP)) A linear map $A$ satisfies the $r$-RIP with constant $\delta_r$, if

$$(1 - \delta_r)\|X\|_2^2 \leq \|A(X)\|_2^2 \leq (1 + \delta_r)\|X\|_2^2$$

is satisfied for all matrices $X \in \mathbb{R}^{m \times r}$ of rank at most $r$.

Theorem 1. Let $X \in \mathbb{R}^{m \times n}$ be a rank $r$ matrix, with singular values $\sigma_1(X) \geq \sigma_2(X) \geq \cdots \geq \sigma_r(X) > 0$ and condition number $\kappa = \sigma_1(X)/\sigma_r(X)$. Denote by $X = A^\dagger \Sigma B$ the corresponding SVD decomposition. Let $U = A^\dagger \Sigma_1^{1/2} \in \mathbb{R}^{m \times r}$ and $V = B^\dagger \Sigma_2^{1/2} \in \mathbb{R}^{n \times r}$. Assume that $A$ satisfies a rank-$6r$ RIP condition with RIP constant $\sigma_{6r} < \frac{1}{\sqrt{m}}$, $\xi_r = \frac{1}{p}$. Then using $T_0 \geq 3\log(\sqrt{Fc}) + 5$ iterations in Algorithm 1 yields a solution $U_0, V_0$ obeying

$$\text{dist}\left(\begin{bmatrix} U_0 \\ V_0 \end{bmatrix}, \begin{bmatrix} U \\ V \end{bmatrix}\right) \leq \frac{1}{4}\sigma_r(U). \quad (8)$$

Furthermore, starting from any initial solution obeying (8), the $i$th iterate of Algorithm 3 satisfies that

$$\text{dist}\left(\begin{bmatrix} U_i \\ V_i \end{bmatrix}, \begin{bmatrix} U \\ V \end{bmatrix}\right) \leq \frac{1}{4}(1 - \gamma_i)\frac{\tilde{\mu}}{\xi - \delta_r}\sigma_r(U). \quad (9)$$

Eq. (8) in the above theorem guarantees that, under the RIP assumption on linear measurements $A$, our algorithm can achieve a good initialization. Eq. (9) states that, starting from a sufficiently accurate initialization, the algorithm exhibits a linear convergence rate. Moreover, the specific convergence rates of the initialization at Eq. (8) and the alternating minimization at Eq. (9) both depend on the RIP constant $\delta_{6r}$.

Jain et al. [14] and Tu et al. [13] built the convergence result for alternating minimization under least square matrix factorization, however, their proving methods cannot be directly adopted to derive the results at Eq. (8). This paper extends [14]’s SVP method to more general objective functions based on M-estimation, and a detailed proof of the convergence of SVP initialization is also provided. For the linear convergence at Eq. (9), Tu et al. [13] gave a similar result, while the gradient descent method was used for the alternating minimization. This paper adopts the proof technique in [17] to establish the linear convergence of the alternating minimization based on Nesterov’s momentum method.

When $\mathcal{M}$ is not smooth, the regularized objective function $\mathcal{M}^\lambda$, defined in Eq. (7), is also not smooth, while the function $\mathcal{M}_p^\lambda$ is smooth. Denote $\{U^+, V^+\} = \min_{U, V} \mathcal{M}^\lambda(U, V)$ and $\{U^{\ast +}, V^{\ast +}\} = \min_{U, V} \mathcal{M}_p^\lambda(U, V)$.

Theorem 2. (Convergence of optimal solution of smoothed objective function) As $\pi \to 0^+$, we have $U^{\pi \dagger}V^{\pi \dagger} \to U^\dagger V^\dagger$.

The above theorem guarantees that the optimal solution of the smoothed object function converges to that of the non-smooth one as the smoothness parameter $\pi$ tends to zero.

Theorem 2 is one of our important contributions. In the literature, most of the state-of-the-art smoothing methods were given without theoretical justifications; see, for example, [18]. Theorem 2 implies that the optimal solution for the smooth approximation indeed converges to the solution of the non-smooth objective function, i.e. the smooth and non-smooth objective functions lead to the same results under some regularity conditions. Under matrix factorization settings, Yang et al. [19] considered Nesterov’s smoothing method to obtain a smooth approximation, while it handled nonnegative matrices only, which is a much simpler problem compared with the one in this paper. Moreover, no strong theoretical convergence guarantee was provided there.

VI. NUMERICAL EXPERIMENTS

As a robust alternative to $L_2$ loss, $L_1$ loss brings computation difficulties since it is not smooth. Using the Nesterov’s smoothing method, we obtained the well-known Huber loss:

$$\mathcal{L}_\mu(a) = \begin{cases} \frac{1}{2\mu}|a|^2 & \text{for } |a| \leq \mu \\ |a| - \frac{\mu}{2} & \text{otherwise} \end{cases},$$

where $\mu$ is the smoothness parameter, and $\mathcal{L}_\mu(a)$ approaches the $L_1$ loss as $\mu$ decreases. In this experiments section, we compare the performance of Huber and $L_2$ loss.

A. Synthetic data

We first generate a matrix $X$ of size $m \times n$ by sampling each entry from the Gaussian distribution $\mathcal{N}(0, 1)$, and the true matrix $X^*$ is obtained via truncated singular value decomposition (SVD) by keeping the first $r$ largest singular values. The sampling matrices $A_i$, $i = 1, \ldots, p$ are independently produced by sampling each entry from the Gaussian distribution $\mathcal{N}(0, 1)$. To evaluate the robustness of the proposed method, we consider eight distributions for the error term $\epsilon$: (a) No error; (b) $\mathcal{N}(0, 2)$; (c) $\mathcal{N}(0, 10)$; (d) $\log\mathcal{N}(0, 1)$; (e) Cauchy; (f) $t(3)$; (g) Pareto$(1, 1)$; (h) $0.9\mathcal{N}(0, 1) + 0.1\mathcal{N}(100, 1)$. Note that (e) and (g) are heavy-tailed with the first-order moment being infinite. We consider two loss functions for the recovery: Huber and $L_2$ loss, and there are 100 replications for each error distribution with $m = n = 100$, $r = 10$ and $p = 5000$.

We consider the following three metrics for evaluation: (1) relative error (RE): $\frac{\|X^* - \hat{U}\hat{V}\|_F}{\|X^*\|_F}$; (2) recovery rate (RR): fraction of the elements whose element-wise relative error is smaller than 5%, where element-wise relative error is defined as $\frac{\|X_{ij} - (\hat{U}\hat{V})_{ij}\|}{\|X_{ij}\|}$; (3) test error (MSE): a test set $\{A^*, b^*\}$ with $p^* = 100$ was generated using the same pipeline, and the test error is defined as $p^*-1\|\mathcal{A}^*(\hat{U}\hat{V}) - b^*\|^2$. For both RE and MSE, a smaller value is desired, while a value closer to one indicates better performance for RR. Table
I presents the results based on the average of 100 replications, and we highlight preferable figures by boldface.

When the data is very heavy-tailed (error $(c)$ and $(g)$), the algorithm diverges for some cases. Table II shows the percentages of convergent repetitions. For the other six distributions, the algorithm converges for all replications. Moreover, we observe that scaling down the learning rate can make the algorithm converge, but this severely slows down the algorithm. When starting with a larger learning rate, the algorithm with $L_2$ loss tends to diverge, while that with the Huber loss converges faster.

![Figure 2: Box plots of recovery rates under Huber (left) and $L_2$ (right) loss functions with eight error distributions: (a) No error; (b) $N(0, 2)$; (c) $N(0, 10)$; (d) log-$N(0, 1)$; (e) Cauchy; (f) $t(3)$; (g) Pareto$(1, 1)$; (h) $0.9 N(0, 1) + 0.1 N(100, 1)$.](image)

Table I: Simulation results for synthetic data with eight error distributions: (a) No error; (b) $N(0, 2)$; (c) $N(0, 10)$; (d) log-$N(0, 1)$; (e) Cauchy; (f) $t(3)$; (g) Pareto$(1, 1)$; (h) $0.9 N(0, 1) + 0.1 N(100, 1)$.

<table>
<thead>
<tr>
<th>Error</th>
<th>Loss</th>
<th>RE</th>
<th>RR</th>
<th>MSE</th>
</tr>
</thead>
<tbody>
<tr>
<td>(a)</td>
<td>Huber</td>
<td>0.003</td>
<td>0.960</td>
<td>0.030</td>
</tr>
<tr>
<td></td>
<td>$L_2$</td>
<td><strong>0.003</strong></td>
<td><strong>0.969</strong></td>
<td><strong>0.022</strong></td>
</tr>
<tr>
<td>(b)</td>
<td>Huber</td>
<td>0.030</td>
<td>0.650</td>
<td>2.885</td>
</tr>
<tr>
<td></td>
<td>$L_2$</td>
<td><strong>0.028</strong></td>
<td><strong>0.671</strong></td>
<td><strong>2.487</strong></td>
</tr>
<tr>
<td>(c)</td>
<td>Huber</td>
<td>0.163</td>
<td>0.185</td>
<td>84.628</td>
</tr>
<tr>
<td></td>
<td>$L_2$</td>
<td><strong>0.141</strong></td>
<td><strong>0.211</strong></td>
<td><strong>63.040</strong></td>
</tr>
<tr>
<td>(d)</td>
<td>Huber</td>
<td>0.028</td>
<td>0.667</td>
<td>2.617</td>
</tr>
<tr>
<td></td>
<td>$L_2$</td>
<td>0.038</td>
<td>0.580</td>
<td>4.571</td>
</tr>
<tr>
<td>(e)</td>
<td>Huber</td>
<td>0.039</td>
<td>0.570</td>
<td>4.832</td>
</tr>
<tr>
<td></td>
<td>$L_2$</td>
<td>1.083</td>
<td>0.037</td>
<td>4742</td>
</tr>
<tr>
<td>(f)</td>
<td>Huber</td>
<td>0.020</td>
<td>0.757</td>
<td>1.238</td>
</tr>
<tr>
<td></td>
<td>$L_2$</td>
<td>0.024</td>
<td>0.711</td>
<td>1.862</td>
</tr>
<tr>
<td>(g)</td>
<td>Huber</td>
<td>0.068</td>
<td>0.399</td>
<td>14.68</td>
</tr>
<tr>
<td></td>
<td>$L_2$</td>
<td>1.825</td>
<td>0.021</td>
<td>12451</td>
</tr>
<tr>
<td>(h)</td>
<td>Huber</td>
<td>0.025</td>
<td>0.704</td>
<td>1.920</td>
</tr>
<tr>
<td></td>
<td>$L_2$</td>
<td>0.494</td>
<td>0.063</td>
<td>791.1</td>
</tr>
</tbody>
</table>

Table II: Percentage of the replications where the algorithm converges for two error distributions: (c) Cauchy; (g) Pareto$(1, 1)$.

<table>
<thead>
<tr>
<th>Error</th>
<th>Huber</th>
<th>$L_2$</th>
</tr>
</thead>
<tbody>
<tr>
<td>(e)</td>
<td>99.4%</td>
<td>80.0%</td>
</tr>
<tr>
<td>(g)</td>
<td>100.0%</td>
<td>62.8%</td>
</tr>
</tbody>
</table>

When there is no error, Huber and $L_2$ loss are comparable, both reaching a recovery rate over 96%. In terms of the metrics presented, $L_2$ loss is better than Huber loss when the error is normally distributed. However, Huber loss is more robust when error introduces bias, skewness or outliers. Huber loss has a substantial advantage over $L_2$ loss especially when the error is skewed (error $(d)$) or heavy-tailed (error $(f)$). When the linear measurements are biased (error $(h)$) or contaminated by heavy-tailed outliers (errors $(e)$ & $(g)$), matrix recovery using $L_2$ recovers only less than 7% of entries and has very large MSE, while employing Huber loss allows the procedure to recover at least 40% of the entries and the recovery error is less than 7% in general.

B. Real data

We further demonstrate the efficiency and robustness of our low-rank matrix recovery algorithm by applying it to two real examples.

The first is a chlorine concentration dataset, available in [20]. The dataset contains a matrix of chlorine concentration levels collected in a water distribution system. The observations in each row are collected at a certain location, and the columns correspond to observations at consecutive time points.

The second is compressed sensing of the MIT logo [14], [21], which is a $38 \times 72$ gray-scale image, see Figure 3.

![Figure 3: MIT logo](image)

1) Chlorine concentration recovery: We use a sub-matrix of size $120 \times 180$ as our ground truth matrix $X^*$ and generated $p = 4200$ sensing matrices and linear measurements according to (2). We set $r = 6$ since the rank-$6$ truncated SVD of $X^*$ achieves a relative low error of 0.069. In this experiment, as in [22], we use outliers to replace measurements $b_i$s, and the outliers are sampled from $N(0, 10\|X^*\|_F)$. The replacement happens with probability 1% or 5%, which is denoted by error distributions $(i)$ and $(j)$, respectively. As a comparison, we
Table III: Performance of Chlorine concentration recovery with four error distributions: (a) No error; (c) $\mathcal{N}(0, 10)$; (i) 1% outliers with outliers sampled from $\mathcal{N}(0, 10)||X^*||_F$; (j) 5% outliers with outliers sampled from $\mathcal{N}(0, 10)||X^*||_F$.

<table>
<thead>
<tr>
<th>Error</th>
<th>Loss</th>
<th>RE</th>
<th>RR</th>
</tr>
</thead>
<tbody>
<tr>
<td>(a)</td>
<td>Huber</td>
<td>0.095</td>
<td>0.130</td>
</tr>
<tr>
<td></td>
<td>$L_2$</td>
<td>0.102</td>
<td>0.121</td>
</tr>
<tr>
<td>(c)</td>
<td>Huber</td>
<td>0.448</td>
<td>0.032</td>
</tr>
<tr>
<td></td>
<td>$L_2$</td>
<td>0.432</td>
<td>0.032</td>
</tr>
<tr>
<td>(i)</td>
<td>Huber</td>
<td>0.097</td>
<td>0.132</td>
</tr>
<tr>
<td></td>
<td>$L_2$</td>
<td>12.587</td>
<td>0.002</td>
</tr>
<tr>
<td>(j)</td>
<td>Huber</td>
<td>0.114</td>
<td>0.109</td>
</tr>
<tr>
<td></td>
<td>$L_2$</td>
<td>N.A.</td>
<td></td>
</tr>
</tbody>
</table>

also consider the cases with no error and the error distribution of $\mathcal{N}(0, 10)$, which correspond to error distributions (a) and (c) in the synthetic data, respectively.

Table III gives the sensing performance under different levels of outliers. The recovery rates are all low, and this possibly is due to the fact that $X^*$ has many entries close to zero. As a result, the relative error might be a more informative metric here. The Huber and $L_2$ loss have comparable performance when there is no error (see also Figure 4) or normal error of a moderate scale (see also Figure 5). Replacing 1% or 5% of the measurements with $\mathcal{N}(0, 10)||X^*||_F$ outliers hardly affects the recovery if Huber loss is used. However, for $L_2$ loss, the result is severely affected by outliers. When there are 5% outliers, the algorithm is either too slow or divergent, thus the result is not available (N.A.).

Figure 4-7 visualize the element-wise comparison between the true matrix and the recovery results by plotting the second row of the matrices in a plot. Figure 4 shows that recovery is good when there is no error. Figure 5 shows that both procedures are affected by the $\mathcal{N}(0, 10)$ noise, but the reconstructed values still share the same pattern with the true values. When there are 1% outliers, using the $L_2$ loss results in large fluctuations and can not recover the matrix (see Figure 6). In contrast, using Huber loss allows us to recover the true $X^*$ even when 5% of the observations are outliers, see Figure 7.

2) Compressed sensing of MIT logo: We use the grayscale logo as the ground truth matrix, which can be well-approximated by a rank-4 matrix (RE = 0.0194 by truncated SVD). We set $m = 38$, $n = 72$, $r = 4$, and take $p = 1200$ measurements using (2). Figure 8 visually compares the sensing results under $L_2$ loss and Huber loss.

If Huber loss is used, the MIT logo can be recovered in all 4 scenarios examined. Recall that the $\mathcal{N}(0, 10)$ noise is the most adverse case for Huber loss among the eight error distributions tested using synthetic dataset. In this MIT logo experiment, the reconstructed image is still recognizable when error follows $\mathcal{N}(0, 10)$. For the other 3 types of errors, the recovery error is hardly noticeable.

On the other hand, using $L_2$ loss allows us to recover the logo when there is no or Normal error. Under Normal error, $L_2$ loss has a slight advantage over Huber loss, as the outcome
The numerical experiments on both synthetic data and real data demonstrate that the Huber loss function provides more robust performance for skewed and/or heavy-tailed data comparing with $L_2$ loss.

Table IV: Recovery of MIT logo with four error distributions: (a) No error; (c) $N(0,10)$; (e) Cauchy; (h) 0.9$N(0,1) + 0.1N(100,1)$.

Appendix

Appendix A: Proof Sketch of Theorem 1

The proof of the convergence shares the same nature with the combined analysis of [13], [23] and [17], thus we only provide a road map.

We first present theorems where $M$ is smooth. The following assumptions are made on the function $M(\cdot)$:

1. Assume that $M(A(Y) - b) - M(A(X) - b) \geq \langle \nabla M(A(X) - b), Y - X \rangle + \xi \|A(X) - A(Y)\|_2^2$. This assumption is the $\xi$-strong convexity assumption in [24].

2. Assume that $M(A(Y) - b) - M(A(X) - b) \leq \langle \nabla M(A(X) - b), Y - X \rangle + \eta \|A(X) - A(Y)\|_2^2$.

3. Without loss of generality, assume $X^*$ is the optimal matrix, then $M(A(X^*) - b) = 0$ and $M(X) \geq 0$.

4. Assume that $M$ also defines a matrix norm when acts on a matrix $X$, and $c_1 \|X\|_2^2 \leq M(X) \leq c_2 \|X\|_2^2$.

Remark: Assumption 1-2 defines the condition number of the function $M(\cdot)$, similar assumptions also appear in [25]. Usually $\kappa = \eta/\xi$ is called the condition number of a function $f$ [26].

Lemma A.1. [27] Let $A$ satisfy $2r$-RIP with constant $\delta_{2r}$. Then for all matrices $X, Y$ of rank at most $r$, we have

$$|\langle A(X), A(Y) \rangle - \langle X, Y \rangle| \leq \delta_{2r} \|X\|_F \|Y\|_F.$$ 

The next lemma characterizes the convergence rate of the initialization procedure:

Lemma A.2. [28] Let $X \in \mathbb{R}^{n \times n}$ be an arbitrary matrix of rank $r$. Let $b = A(X) \in \mathbb{R}^p$ be $p$ linear measurements. Consider the iterative updates

$$Y^{t+1} \leftarrow \mathcal{P}_r \left( Y^t - \xi_t \nabla_X M(A(Y^t) - b) \right),$$

where $Y^t$ are $m \times n$ matrices. Then

$$\|Y^t - X\|_F \leq \psi(A)^t \|Y^0 - X\|_F.$$
holds, where $\psi(A)$ is defined as
\[
\psi(A) = 2\sup_x \|x\|_F = \|Y\|_F = 1, \text{rank}(X) \leq 2r, \text{rank}(Y) \leq 2r
\]
\[
|x, A(Y) - (X, Y)|.
\]

We can prove this lemma by using the results of Theorem A.1.

First, we prove that the initialization procedure indeed converges to the true value of $X$.

**A. Proof of Eq. (8) in Theorem 1**

**Lemma A.3** (Initialization). Assume that $\xi \frac{1 + \delta_{2k}}{2 \delta_{2k}} - \eta > 0$.

Denote $\tilde{M}(X) = M(A(X) - b)$. Let $X^*$ be an optimal solution and let $X^t$ be the value obtained by Algorithm 1 at $t^{th}$ iteration. Then

\[
\tilde{M}(X^{t+1}) \leq \tilde{M}(X^t) + \left(\frac{1 + \delta_{2k}}{1 - \delta_{2k}} - \eta \right) \|A(X^t) - X^t\|_F^2.
\]

**Proof.** From assumption, we have

\[
\tilde{M}(X^{t+1}) - \tilde{M}(X^t)
\]
\[
\leq \langle \nabla \tilde{M}(X^t), X^{t+1} - X^t \rangle + \|\tilde{M}(X^{t+1}) - \tilde{M}(X^t)\|_F^2
\]
\[
\leq \langle \nabla \tilde{M}(X^t), X^{t+1} - X^t \rangle + \left(1 + \delta_{2k}\right) \|X^{t+1} - X^t\|_F^2
\]
\[
\text{where the last inequality comes from RIP. Let } Y^{t+1} = X^t - \frac{1}{2\xi(1 + \delta_{2k})} \nabla \tilde{M}(X^t), \text{ and }
\]
\[
f_t(X) = \langle \nabla \tilde{M}(X^t), X - X^t \rangle + \left(1 + \delta_{2k}\right) \|X - X^t\|_F^2.
\]

Then
\[
f_t(X) = \xi \left(1 + \delta_{2k}\right) \left[\|X - Y^{t+1}\|_F^2 - \frac{1}{4\xi^2(1 + \delta_{2k})^2} \left\|\nabla \tilde{M}(X^t)\right\|_F^2\right]
\]

By definition, $P_k(Y^{t+1}) = X^{t+1}$, then $f_t(X^{t+1}) \leq f_t(X^*)$.

Thus
\[
\tilde{M}(X^{t+1}) - \tilde{M}(X^t) \leq f_t(X^{t+1}) \leq f_t(X^*)
\]
\[
= \langle \nabla \tilde{M}(X^t), X^* - X^t \rangle + \left(1 + \delta_{2k}\right) \|X^* - X^t\|_F^2
\]
\[
\leq \nabla \tilde{M}(X^t), X^* - X^t \rangle + \left(1 + \delta_{2k}\right) \|A(X^t) - A(X^*)\|_F^2
\]
\[
\leq \left(\frac{1 + \delta_{2k}}{1 - \delta_{2k}} - \eta \right) \|A(X^*) - A(X^t)\|_F^2
\]
\[
+ \left(\frac{1 + \delta_{2k}}{1 - \delta_{2k}} - \eta \right) \|A(X^*) - A(X^t)\|_F^2.
\]

**Theorem A.1.** Let $b = A(X^*) + e$ for rank $k$ matrix $X^*$ for an error vector $e \in \mathbb{R}^p$, $D = \frac{C^2}{2} + \left(\frac{1 + \delta_{2k}}{1 - \delta_{2k}} - \eta \right) \left(\frac{C^2}{2} + \sqrt{\frac{2}{C_0^2} + \frac{1}{C_1}}\right)$, Then, under the assumption that $D < 1$, Algorithm I with step size $\eta_t = \frac{C^2(1 + \delta_{2k})}{2\xi(1 + \delta_{2k})}$ outputs a matrix $X$ of rank at most $k$ such that $\tilde{M}(A(X^t) - b) \leq \frac{1}{C^2 + \varepsilon} \left[\frac{C^2 + \varepsilon}{2\xi(1 + \delta_{2k})}\right]^2$, where $\varepsilon > 0$, in at most $\frac{1}{\log D} \log \frac{(C^2 + \varepsilon)[\varepsilon]^2}{2\xi(1 + \delta_{2k})}$ iterations.

**Proof.** Let the current solution $X^t$ satisfy $M(X^t) \geq \frac{C^2}{2} \|e\|^2$.

By lemma A.3 and $b - A(X^*) = e$, we have
\[
M(X^{t+1}) \leq \frac{\|e\|^2}{2} + \left(1 + \frac{\delta_{2k}}{1 - \delta_{2k}} - \eta \right) \|b - A(X^t) - e\|^2
\]
\[
\leq \frac{\|e\|^2}{2} + \left(1 + \frac{\delta_{2k}}{1 - \delta_{2k}} - \eta \right) \|b - A(X^t)\|^2
\]
\[
- \frac{\|e\|^2}{2} \leq \frac{\|M(X^t)\|}{C^2} + \left(1 + \frac{\delta_{2k}}{1 - \delta_{2k}} - \eta \right) \left(\frac{2M(X^t)}{C^2}\right)
\]
\[
+ \frac{M(X^t)}{c_1} + \frac{\sqrt{2M(X^t)}}{C\sqrt{c_1}} = DM(X^t).
\]

Since $D < 1$, combining the fact that $M(X^0) \leq c_2 \|b\|^2$, by taking $t = \left[\frac{1}{\log D} \log \frac{(C^2 + \varepsilon)[\varepsilon]^2}{2\xi(1 + \delta_{2k})}\right]$, we complete the proof. □

The following lemma is adapted from [13]:

**Lemma A.4.** Let $X_1, X_2 \in \mathbb{R}^{m \times n}$ be two rank $r$ matrices with SVD decomposition $X_1 = A_1 \Sigma_1 B_1$, $X_2 = A_2 \Sigma_2 B_2$.

For $l = 1, 2$, define $U_l = A_l^\top \Sigma_l^{1/2} \in \mathbb{R}^{m \times r}$, $V_l = B_l^\top \Sigma_l^{1/2} \in \mathbb{R}^{r \times n}$.

Assume $X_1, X_2$ obey $\|X_2 - X_1\|_F \leq \frac{1}{4} \sigma_r(X_1)$. Then
\[
dist^2 \left(\begin{bmatrix} U_2 \\ V_2 \end{bmatrix}, \begin{bmatrix} U_1 \\ V_1 \end{bmatrix}\right) \leq \frac{2}{\sqrt{2} - 1} \frac{\|X_2 - X_1\|_F^2}{\sigma_r(X_1)}.
\]

Combine Lemma A.1 and A.4, and follow a similar route of [13], we can prove that using more than $3 \log(\sqrt{r} \kappa) + 5$ iterations of Algorithm 1, we can obtain
\[
dist \left(\begin{bmatrix} U_0 \\ V_0 \end{bmatrix}, \begin{bmatrix} U \\ V \end{bmatrix}\right) \leq \frac{1}{4} \sigma_r(U).
\]

Thus we finished the proof of convergence in the initialization procedure. Next, we prove the linear convergence of the main algorithm on the penalized objective function.

We first rewrite the objective function (7) by using the uplifting technique, so that it is easier to simultaneously consider $M$ and the regularization term. To see this, consider rank $r$ matrix $X \in \mathbb{R}^{m \times n}$ with SVD decomposition $X = U^\top \Sigma V$.

Define $\text{Sym} : \mathbb{R}^{m \times n} \to \mathbb{R}^{(m+n) \times (m+n)}$ as
\[
\text{Sym}(X) = \begin{bmatrix} 0_{m \times m} & X \\ X^\top & 0_{n \times n} \end{bmatrix}.
\]

Given the block matrix $A = \begin{bmatrix} A_{11} & A_{12} \\ A_{21} & A_{22} \end{bmatrix}$ with $A_{11} \in \mathbb{R}^{m \times m}$, $A_{12} \in \mathbb{R}^{m \times n}$, $A_{21} \in \mathbb{R}^{n \times m}$, $A_{22} \in \mathbb{R}^{n \times n}$, define $\text{P}_{\text{diag}}(A) = \begin{bmatrix} A_{11} & 0_{m \times n} \\ 0_{n \times m} & A_{22} \end{bmatrix}$ and $\text{P}_{\text{off}}(A) = \begin{bmatrix} 0_{m \times m} & A_{12} \\ A_{21} & 0_{n \times n} \end{bmatrix}$.

Define $B : \mathbb{R}^{(m+n) \times (m+n)} \to \mathbb{R}^p$ as the uplifted version of the $A$ operator
\[
B(X)_k = \langle b_k, X \rangle, \text{ where } b_k = \text{Sym}(A_k).
\]
Define $W = [U^T, V^T]$. As a result, we can rewrite the objective function (7) as:

$$g(W) = g(U, V) = M(b - A(U^T V)) + \lambda \|UU^T - VV^T\|_F^2 = \frac{1}{2} M(B(\text{Sym}(U^T V)) - \text{Sym}(X)) + \frac{1}{2} \lambda \|\text{Sym}(U^T V) - \text{Sym}(X)\|_F^2.$$  

From equation (10) we can see that two parts of the penalized objective function have similar structures.

As a result, although we made Assumptions 1, 2, 4 on the original function $M$, we can see from (10) that the penalized objective function still retains a similar property. As for Assumption 3, we can also use some location transform techniques to make the penalized objective function satisfies this assumption. Thus, we can deal with the penalized objective function or the unpenalized objective function in a similar way.

Then, similar to [13], the alternating minimization together with Nesterov’s momentum algorithm with respect to $U$ and $V$ (sub-matrices of $W$) can be written as a NAG algorithm applied to $g(W)$ with respect to $W$.

For the convergence analysis of the Nesterov’s momentum algorithm, we employ the following lemma, which is a theorem in [17].

**Lemma A.5.** For minimization problem $\min_{x \in \mathcal{X}} f(x)$, where $x$ is a vector, using the Lyapunov function

$$\tilde{V}_k = f(y_k) + \xi \|z_k - x^*\|^2,$$

it can be shown that

$$\tilde{V}_{k+1} - \tilde{V}_k = -\tau_k \tilde{V}_k + \varepsilon_{k+1}$$  \hspace{1cm} (11)

where the error is expressed as

$$\varepsilon_{k+1} = \left(\frac{\tau_k^2}{4\xi}\right) \|\nabla f(x_k)\|^2 + \left(\tau_k \eta - \frac{\xi}{\tau_k}\right) \|x_k - y_k\|^2,$$

$\tau_k$ is the step size in Nesterov’s momentum algorithm, usually equals $1/\sqrt{\kappa}$, $y_{k+1} = x_k - \frac{1}{\kappa} \nabla f(x_k)$, $x_{k+1} = \frac{1}{1+\tau_k} y_k + \frac{\tau_k}{1+\tau_k} z_k$, $z_{k+1} = z_k + \tau_k \left(x_{k+1} - z_k - \frac{1}{2\xi} \nabla f(x+k+1)\right)$. Assume that $\tau_0 = 0$, $\tau_1 = \tau_2 = \cdots = \tilde{\tau}$, and $\varepsilon_1, \cdots, \varepsilon_{k+1}$ has a common upper bound $\tilde{\varepsilon}$, then (11) implies:

$$|\tilde{V}_{k+1}| = |(1 - \tilde{\tau})^{k+1} \tilde{V}_0 + \sum_{i=1}^{k+1} (1 - \tilde{\tau})^{i-1} \varepsilon_{k+2-i}| \leq (1 - \tilde{\tau})^{k+1} |\tilde{V}_0| + \frac{\tilde{\epsilon} - \tilde{\varepsilon}(1 - \tilde{\tau})}{\tilde{\tau}}.$$

Substitute $x_k$ with $W_{k+1} = [U^{k+1}, V^{k+1}]^T$, $f$ with $g$. To deal with the convergence analysis with respect to $W_0 - W^*$ and $W_0 - W^*$, we need to handle two parts. The first part is the error part with respect to $\tilde{\varepsilon}$. This can be solved by choosing an initial estimate close to the true value. As a result, $\|\nabla f(x_0)\|$ can be arbitrary close to 0. For notational simplicity, assume that $\frac{\tilde{\varepsilon} - \tilde{\varepsilon}(1 - \tilde{\tau})}{\tilde{\tau}} \leq \varepsilon^\dagger$. Since $\tilde{V}_k$ still satisfies Assumption 1 and 2, without loss of generality, assume the corresponding parameter are $\tilde{\xi}$ and $\tilde{\eta}$.

Next we want to study the relation between $W_i - W^*$ and $\tilde{V}_i$. This involves Assumption 1 and 2 as well as Lemma A.1. With rough handling of the gradient part in Assumption 1 and 2, we can obtain

$$\tilde{\xi} \|A(W_i) - A(W^*)\|_2^2 \leq (1 - \tilde{\tau})^i \tilde{\mu} \|A(W_0) - A(W^*)\|_2^2 + \tilde{\varepsilon}.$$  

Notice that $\tilde{\varepsilon}$ can be made arbitrary small so that

$$\tilde{\xi} \|A(W_i) - A(W^*)\|_2^2 \leq (1 - \tilde{\tau})^i \tilde{\mu} \|A(W_0) - A(W^*)\|_2^2$$  

and $1 - \tilde{\tau}$ can still be larger than 0 smaller than 1. Employ the RIP property, we have

$$\tilde{\xi}(1 - \delta_r) \|W_i - W^*\|_2^2 \leq (1 - \tilde{\tau})^i \tilde{\mu} (1 + \delta_r) \|W_0 - W^*\|_2^2.$$  

Thus

$$\text{dist} \left( \left[ \begin{array}{c} U_i \\ V_i \end{array} \right], \left[ \begin{array}{c} U \\ V \end{array} \right] \right) \leq (1 - \tilde{\tau})^i \tilde{\mu} \left( \frac{1 + \delta_r}{\xi - \delta_r} \right) \text{dist} \left( \left[ \begin{array}{c} U_0 \\ V_0 \end{array} \right], \left[ \begin{array}{c} U \\ V \end{array} \right] \right) \leq \frac{1}{4} (1 - \tilde{\tau})^i \tilde{\mu} \left( \frac{1 + \delta_r}{\xi - \delta_r} \sigma_r(U). \right)$$

Theorem 1 is proved.

**Remark on (10):** From (10), we provide a guideline with respect to the selection of $\lambda$ compared with [13], by combining (10) and Assumption 4.

**Appendix B: Proof Sketch of Theorem 2**

**Proof.** Employ Theorem 1 in [29] and Lemma 3 in [30] we know that

$$M^\lambda(b - A(U^{\pi \ast T}V^{\pi \ast})) \rightarrow M^\lambda(b - A(U^{\ast T}V^{\ast}))$$

Giving the fact that $M^\lambda$ and $M^\lambda$ are convex, as well as the restricted isometry property, we finish the proof of Theorem 2. \hfill \Box

**REFERENCES**


