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# GRADIENT BASED BLOCK COORDINATE DESCENT ALGORITHMS FOR JOINT APPROXIMATE DIAGONALIZATION OF MATRICES 

JIANZE LI ${ }^{\dagger}$, KONSTANTIN USEVICH ${ }^{\ddagger}$, AND PIERRE COMON ${ }^{\natural}$


#### Abstract

In this paper, we propose a gradient based block coordinate descent (BCD-G) framework to solve the joint approximate diagonalization of matrices defined on the product of the complex Stiefel manifold and the special linear group. Instead of the cyclic fashion, we choose the block for optimization in a way based on the Riemannian gradient. To update the first block variable in the complex Stiefel manifold, we use the well-known line search descent method. To update the second block variable in the special linear group, based on four different kinds of elementary transformations, we construct three classes: GLU, GQU and GU, and then get three BCD-G algorithms: BCD-GLU, BCD-GQU and BCD-GU. We establish the global convergence and weak convergence of these three algorithms using the Lojasiewicz gradient inequality under the assumption that the iterates are bounded. We also propose a gradient based Jacobi-type framework to solve the joint approximate diagonalization of matrices defined on the special linear group. Similar as in the BCD-G case, using the GLU and GQU classes of elementary transformations, we focus on the Jacobi-GLU and JacobiGQU algorithms, and establish their global convergence and weak convergence as well. All the algorithms and convergence results in this paper also apply to the real case.


## 1. Introduction

1.1. Problem formulation. Let $1 \leq m \leq n$. For a matrix $\boldsymbol{X} \in \mathbb{C}^{n \times m}$, we denote by $\boldsymbol{X}^{\boldsymbol{\top}}$, $\boldsymbol{X}^{*}$ and $\boldsymbol{X}^{\mathrm{H}}$ its transpose, conjugate and conjugate transpose, respectively. We shall also use $(\cdot)^{\star}$ to denote either $(\cdot)^{\mathrm{T}}$ or $(\cdot)^{\mathrm{H}}$. Let $\left\{\boldsymbol{A}^{(\ell)}\right\}_{1 \leq \ell \leq \mathrm{L}} \subseteq \mathbb{C}^{n \times m}$ be a set of complex matrices, and $\mu_{\ell} \in \mathbb{R}^{+}$for $1 \leq \ell \leq L$. It is well-known that the blind source separation (BSS) [16, 17, 35, 42] can be formulated as finding a full column rank matrix $Z \in \mathbb{C}^{n \times m}$ to make the matrices $\boldsymbol{W}^{(\ell)}=\boldsymbol{Z}^{\star} \boldsymbol{A}^{(\ell)} \boldsymbol{Z} \in \mathbb{C}^{m \times m}$ simultaneously as diagonal as possible. This is to solve the joint approximate diagonalization of matrices (JADM) problem, which consists in minimizing the cost function

$$
\begin{equation*}
f(\boldsymbol{Z})=\sum_{\ell=1}^{L} \mu_{\ell}\left\|\operatorname{offdiag}\left\{\boldsymbol{W}^{(\ell)}\right\}\right\|^{2}, \tag{1}
\end{equation*}
$$

[^0]where offdiag $\{\cdot\}$ is the zero diagonal operator, which sets all the diagonal elements of a square matrix in $\mathbb{C}^{m \times m}$ to zero.

Note that, if $\boldsymbol{Z} \in \mathbb{C}^{n \times m}$ is of full column rank and $\lambda \in \mathbb{C}$ is nonzero, then $\lambda \boldsymbol{Z}$ is still of full column rank and the function $f$ in (1) satisfies $f(\lambda \boldsymbol{Z})=|\lambda|^{2} f(\boldsymbol{Z})$. Therefore, the optimization problem of minimizing $f$ on the set of full column rank matrices is not well defined. To tackle this issue, one approach is to introduce new cost functions. For example, several cost functions with the scale and permutation invariance property are introduced in [4, 48]. In this paper, we still use the function (1), while restricting the variable $\boldsymbol{Z}$ in a smaller set $\mathbf{R S L}(m, n, \mathbb{C})$, which will be defined in (2).

Let $\mathbf{G L}_{m}(\mathbb{C}) \stackrel{\text { def }}{=}\left\{\boldsymbol{X} \in \mathbb{C}^{m \times m}, \operatorname{det}(\boldsymbol{X}) \neq 0\right\}$ be the general linear group, and $\mathbf{S L}_{m}(\mathbb{C}) \stackrel{\text { def }}{=}$ $\left\{\boldsymbol{X} \in \mathbf{G L}_{m}(\mathbb{C}), \operatorname{det}(\boldsymbol{X})=1\right\}$ be the special linear group. We define the rectangular special linear set as

$$
\begin{equation*}
\mathbf{R S L}(m, n, \mathbb{C}) \stackrel{\text { def }}{=}\left\{\boldsymbol{X} \in \mathbb{C}^{n \times m}, \boldsymbol{X}^{\mathrm{H}} \boldsymbol{X} \in \mathbf{S L}_{m}(\mathbb{C})\right\} \tag{2}
\end{equation*}
$$

In this paper, we mainly study the JADM problem on $\mathbf{R S L}(m, n, \mathbb{C})$, which is to minimize the cost function

$$
\begin{equation*}
f(\boldsymbol{Z})=\sum_{\ell=1}^{L} \mu_{\ell}\left\|\operatorname{offdiag}\left\{\boldsymbol{W}^{(\ell)}\right\}\right\|^{2} \tag{3}
\end{equation*}
$$

where $\boldsymbol{Z} \in \mathbf{R S L}(m, n, \mathbb{C})$.

### 1.2. Two equivalent reformulations.

1.2.1. First reformulation on $\mathbf{S t}(m, n, \mathbb{C}) \times \mathbf{S L}_{m}(\mathbb{C})$. Let $\mathbf{S t}(m, n, \mathbb{C}) \stackrel{\text { def }}{=}\left\{\boldsymbol{Y} \in \mathbb{C}^{n \times m}, \boldsymbol{Y}^{\mathrm{H}} \boldsymbol{Y}=\right.$ $\left.\boldsymbol{I}_{m}\right\}$ be the complex Stiefel manifold. We have the following relationship ${ }^{1}$ :

$$
\begin{equation*}
\mathbf{R S L}(m, n, \mathbb{C})=\mathbf{S t}(m, n, \mathbb{C}) \times \mathbf{S L}_{m}(\mathbb{C}) \tag{4}
\end{equation*}
$$

By equation (4), problem (3) is equivalent to minimizing the cost function

$$
\begin{equation*}
f: \mathbf{S t}(m, n, \mathbb{C}) \times \mathbf{S L}_{m}(\mathbb{C}) \rightarrow \mathbb{R}^{+}, \quad(\boldsymbol{Y}, \boldsymbol{X}) \mapsto \sum_{\ell=1}^{L} \mu_{\ell}\left\|\operatorname{offdiag}\left\{\boldsymbol{W}^{(\ell)}\right\}\right\|^{2} \tag{5}
\end{equation*}
$$

where $\boldsymbol{W}^{(\ell)}=(\boldsymbol{Y} \boldsymbol{X}) \boldsymbol{A}^{(\ell)}(\boldsymbol{Y} \boldsymbol{X}) \in \mathbb{C}^{m \times m}$.
1.2.2. Second reformulation on $\mathbf{S L}_{m}(\mathbb{C})$. In problem (5), it is often possible to find a solution $\boldsymbol{Y}_{*} \in \mathbf{S t}(m, n, \mathbb{C})$ in advance by some other method, e.g. PCA $[15,16,17]$. This well-known procedure allows to reduce the space dimension by projection onto the dominant subspace, and at the same time to reduce the noise level. Let $\boldsymbol{B}^{(\ell)}=\boldsymbol{Y}_{*}^{\star} \boldsymbol{A}^{(\ell)} \boldsymbol{Y}_{*}$ for $1 \leq \ell \leq L$. Then the cost function (5) becomes

$$
\begin{equation*}
g: \mathbf{S L}_{m}(\mathbb{C}) \rightarrow \mathbb{R}^{+}, \quad \boldsymbol{X} \mapsto g(\boldsymbol{X})=\sum_{\ell=1}^{L} \mu_{\ell}\left\|\operatorname{offdiag}\left\{\boldsymbol{W}^{(\ell)}\right\}\right\|^{2} \tag{6}
\end{equation*}
$$

[^1]where $\boldsymbol{W}^{(\ell)}=\boldsymbol{X}^{\diamond} \boldsymbol{B}^{(\ell)} \boldsymbol{X} \in \mathbb{C}^{m \times m}$. A complex matrix $\boldsymbol{B} \in \mathbb{C}^{m \times m}$ is called Hermitian if $\boldsymbol{B}^{\mathrm{H}}=\boldsymbol{B}$. It is called complex symmetric if $\boldsymbol{B}^{\top}=\boldsymbol{B}$. Problem (6) has the following well-known special cases:

- joint approximate diagonalization of Hermitian matrices (JADM-H)[41, 35]: $(\cdot)^{\wedge}=(\cdot)^{\mathrm{H}}$, $\boldsymbol{B}^{(\ell)}$ is Hermitian and $\mu_{\ell}=1$ for $1 \leq \ell \leq L$;
- joint approximate diagonalization of complex symmetric matrices (JADM-CS) [35]: $(\cdot)^{\uparrow}=(\cdot)^{\top}, \boldsymbol{B}^{(\ell)}$ is complex symmetric and $\mu_{\ell}=1$ for $1 \leq \ell \leq L$;
- joint approximate diagonalization of real symmetric matrices (JADM-RS) [4, 5]: on the real field $\mathbb{R},(\cdot)^{\star}=(\cdot)^{\top}, \boldsymbol{B}^{(\ell)}$ is real symmetric and $\mu_{\ell}=1$ for $1 \leq \ell \leq L$.
These three problems have been widely used in the BSS and Independent component analysis (ICA) [13, 17, 4, 6]. Compared with the joint approximate diagonalization by orthogonal transformations $[11,12,16,24,27,25,26,44]$, the non-orthogonal joint diagonalizer in (6) does not require prewhitening and may less suffer from noise. However, since $\mathbf{S L}_{m}(\mathbb{C})$ is not compact, solving this problem is much more difficult.

To solve the JADM-RS problem, Jacobi-type algorithms were introduced based on the LU and QR decompositions in [5], and based on the Givens transformations, hyperbolic transformations and diagonal transformations in [42, Eq. (9)]. To solve the JADM-H problem, Jacobi-type algorithms were proposed based on the LU decomposition in [35, 36], and based on the QL decomposition in [41]. To solve the JADM-CS problem, a Jacobi-type algorithm was proposed based on the LU decomposition in [34, 35]. However, to our knowledge, there was no theoretical result about the convergence of these Jacobi-type algorithms in the literature.
1.3. Block coordinate descent. Suppose that $\left\{\mathcal{M}_{i}\right\}_{1 \leq i \leq d}$ is a set of smooth manifolds. To minimize the following smooth function

$$
\begin{equation*}
f: \mathcal{M}_{1} \times \mathcal{M}_{2} \times \cdots \times \mathcal{M}_{d} \longrightarrow \mathbb{R}^{+} \tag{7}
\end{equation*}
$$

a popular approach is the block coordinate descent (BCD) [8, 31, 32, 46, 47, 28]. In this method, only one block variable is updated at each iteration, while other block variables are fixed. Then problem (7) is decomposed to a sequence of lower-dimensional optimization problems, similar as the subproblems in Jacobi-type algorithms. In BCD algorithm, there are different ways to choose blocks for optimization, including the essentially cyclic, cyclic, random fashions [46, 47] and the so-called maximum block improvement (MBI) method [14, 29].
1.4. BCD-G algorithm for the first reformulation on $\mathbf{S t}(m, n, \mathbb{C}) \times \mathbf{S L}_{m}(\mathbb{C})$. Suppose that

$$
\begin{equation*}
f: \mathbf{S t}(m, n, \mathbb{C}) \times \mathbf{S L}_{m}(\mathbb{C}) \rightarrow \mathbb{R}^{+} \tag{8}
\end{equation*}
$$

is an abstract smooth function, which includes the cost function (5) as a special case. For $\boldsymbol{\omega}=(\boldsymbol{Y}, \boldsymbol{X}) \in \mathbf{S t}(m, n, \mathbb{C}) \times \mathbf{S L}_{m}(\mathbb{C})$, we define

$$
\begin{equation*}
f_{1, \boldsymbol{X}}: \boldsymbol{Y} \mapsto f(\boldsymbol{Y}, \boldsymbol{X}), \quad f_{2, \boldsymbol{Y}}: \boldsymbol{X} \mapsto f(\boldsymbol{Y}, \boldsymbol{X}) \tag{9}
\end{equation*}
$$

as the restricted functions. For simplicity, we denote the Riemannian $\operatorname{gradient}^{2} \operatorname{grad} f_{1}(\boldsymbol{\omega}) \stackrel{\text { def }}{=}$ $\operatorname{grad} f_{1, \boldsymbol{X}}(\boldsymbol{Y})$ and $\operatorname{grad} f_{2}(\boldsymbol{\omega}) \stackrel{\text { def }}{=} \operatorname{grad} f_{2, \boldsymbol{Y}}(\boldsymbol{X})$. Let $\boldsymbol{\omega}_{k} \stackrel{\text { def }}{=}\left(\boldsymbol{Y}_{k}, \boldsymbol{X}_{k}\right)$ for $k \geq 0$. Now we propose

[^2]the following gradient based block coordinate descent (BCD-G) algorithm to minimize the cost function (8). It is easy to see that, in Algorithm 1, we can always choose $t_{k}=1$ or 2 such that the inequality (10) is satisfied ${ }^{3}$.

```
Algorithm 1: BCD-G algorithm
    Input: A starting point \(\boldsymbol{\omega}_{0}, 0<v<\frac{\sqrt{2}}{2}\).
    Output: Sequence of iterates \(\boldsymbol{\omega}_{k}\).
    for \(k=1,2, \cdots\), do
        Choose \(t_{k}=1\) or 2 such that
                                    \(\left\|\operatorname{grad} f_{t_{k}}\left(\boldsymbol{\omega}_{k-1}\right)\right\| \geq v\left\|\operatorname{grad} f\left(\boldsymbol{\omega}_{k-1}\right)\right\| ;\)
        if \(t_{k}=1\) then
            Update \(\boldsymbol{U}_{k}\) using the line search descent method (cf. Equation (31) in Section 3);
            Set \(\boldsymbol{X}_{k}=\boldsymbol{X}_{k-1}\);
        else
            Set \(\boldsymbol{U}_{k}=\boldsymbol{U}_{k-1}\);
            Update \(\boldsymbol{X}_{k}\) using the elementary transformations (cf. Subalgorithm 1a to 1c in
            Section 4);
        end if
    end for
```

Now we need to recall some Lie groups before introducing further details about Algorithm 1. A matrix $\boldsymbol{X} \in \mathbb{C}^{m \times m}$ is said to be upper triangular if $\mathrm{X}_{i j}=0$ for $i>j$. Let $\mathbf{U T}_{m}(\mathbb{C}) \subseteq \mathbf{G L}_{m}(\mathbb{C})$ be the upper triangular subgroup, that is

$$
\mathbf{U T}_{m}(\mathbb{C}) \stackrel{\text { def }}{=}\left\{\boldsymbol{X} \in \mathbf{G L}_{m}(\mathbb{C}), \boldsymbol{X} \text { is upper triangular }\right\} .
$$

Let $\mathbf{E U T}_{m}(\mathbb{C})=\mathbf{U T}_{m}(\mathbb{C}) \cap \mathbf{S L}_{m}(\mathbb{C})$, i.e., the set of upper triangular matrices with determinant equal to 1 . Similarly, we let $\mathbf{L T} T_{m}(\mathbb{C}) \subseteq \mathbb{C}^{m \times m}$ be the lower triangular subgroup and $\mathbf{E L T}_{m}(\mathbb{C})=$ $\mathbf{L T}_{m}(\mathbb{C}) \cap \mathbf{S L}_{m}(\mathbb{C})$. Let $\mathbf{U}_{m} \subseteq \mathbb{C}^{m \times m}$ be the unitary group, and $\mathbf{S U}_{m} \subseteq \mathbf{U}_{m}$ be the special unitary group.

In Algorithm 1, to update $\boldsymbol{U}_{k}$, we choose the line search descent method [3], which will be presented in Section 3. To update $\boldsymbol{X}_{k}$, using the Givens plane, plane upper triangular, plane lower triangular and plane diagonal transformations ${ }^{4}$, which will be introduced in Section 4, we construct three classes of elementary transformations: GLU, GQU and GU.

- Any matrix $\boldsymbol{X} \in \mathbf{S L}_{m}(\mathbb{C})$ has the $L U$ decomposition $[18] \boldsymbol{X}=\boldsymbol{L} \boldsymbol{U}$ with $\boldsymbol{L} \in \mathbf{L T}_{m}(\mathbb{C})$ and $\boldsymbol{U} \in \mathbf{U T}_{m}(\mathbb{C})$. The first class GLU is inspired by the following relationship

$$
\begin{equation*}
\mathbf{R S L}(m, n, \mathbb{C})=\mathbf{S t}(m, n, \mathbb{C}) \times \mathbf{E L T}_{m}(\mathbb{C}) \times \mathbf{E U T}_{m}(\mathbb{C}) \tag{11}
\end{equation*}
$$

[^3]which follows from equation (4) and the LU decomposition, and it includes the plane lower triangular, plane upper triangular and plane diagonal transformations (Subalgorithm 1a). In this case, we call Algorithm 1 the $B C D-G L U$ algorithm.

- Any matrix $\boldsymbol{X} \in \mathbf{S L}_{m}(\mathbb{C})$ has the the $Q U$ decomposition $\boldsymbol{X}=\boldsymbol{Q U}$ with $\boldsymbol{Q} \in \mathbf{S U}_{m}$ and $\boldsymbol{U} \in \mathbf{U T}_{m}(\mathbb{C})$. The second class GQU is inspired by the following relationship

$$
\begin{equation*}
\mathbf{R S L}(m, n, \mathbb{C})=\mathbf{S t}(m, n, \mathbb{C}) \times \mathbf{S U}_{m} \times \mathbf{E U T}_{m}(\mathbb{C}) \tag{12}
\end{equation*}
$$

which follows from equation (4) and the QU decomposition, and it includes the Givens plane, plane upper triangular and plane diagonal transformations (Subalgorithm 1b). In this case, we call Algorithm 1 the $B C D-G Q U$ algorithm.

- Note that the equation (12) can be further written as

$$
\begin{equation*}
\mathbf{R S L}(m, n, \mathbb{C})=\mathbf{S t}(m, n, \mathbb{C}) \times \mathbf{S U}_{m} \times \mathbf{E U T}_{m}(\mathbb{C})=\mathbf{S t}(m, n, \mathbb{C}) \times \mathbf{E U T}_{m}(\mathbb{C}) \tag{13}
\end{equation*}
$$

We construct the third class GU, which only includes the plane upper triangular and plane diagonal transformations (Subalgorithm 1c). In this case, we call Algorithm 1 the $B C D-G U$ algorithm.

### 1.5. Jacobi-G algorithm for the second reformulation on $\mathbf{S L}_{m}(\mathbb{C})$. Suppose that

$$
\begin{equation*}
g: \mathbf{S L}_{m}(\mathbb{C}) \rightarrow \mathbb{R}^{+} \tag{14}
\end{equation*}
$$

is an abstract smooth function, which includes the cost function (6) as a special case. To minimize the cost function (14), using the Givens plane, plane upper triangular, plane lower triangular and plane diagonal transformations, which will be introduced in Section 4, we will propose two types of gradient-based Jacobi-type algorithms: Jacobi-GLU and Jacobi-GQU, which will be detailedly formulated in Section 7.3.1. Now we only show the motivations of these two algorithms. Jacobi-GLU is inspired by the LU decomposition

$$
\begin{equation*}
\mathbf{S L}_{m}(\mathbb{C})=\mathbf{E L T}_{m}(\mathbb{C}) \times \mathbf{E U T}_{m}(\mathbb{C}) \tag{15}
\end{equation*}
$$

and it uses the GLU class. Jacobi-GQU is inspired by the QU decomposition

$$
\begin{equation*}
\mathbf{S L}_{m}(\mathbb{C})=\mathbf{S U}_{m} \times \mathbf{E U T}_{m}(\mathbb{C}) \tag{16}
\end{equation*}
$$

and it uses the GQU class.
1.6. Contributions. The main contributions of this paper can be summarized as follows:

- To solve problem (5), we propose the BCD-G framework (Algorithm 1), which chooses the block for optimization in a way based on the Riemannian gradient. This is similar to the gradient-based way of choosing index pairs in Jacobi-G algorithms [20, 24, 44].
- To update the first block variable $\boldsymbol{Y}$, we adopt the well-known line search descent method. To update the second block variable $\boldsymbol{X}$, based on four kinds of elementary transformations, we construct three classes (GLU, GQU and GU), and then get three BCD-G algorithms: BCD-GLU, BCD-GQU and BCD-GU.
- We establish the global convergence ${ }^{5}$ and weak convergence ${ }^{6}$ of BCD-GLU, BCD-GQU and BCD-GU algorithms using the Eojasiewicz gradient inequality, under the assumption that the iterates are bounded, that is, there exists $\mathrm{M}_{\omega}>0$ such that

$$
\begin{equation*}
\left\|\boldsymbol{\omega}_{k}\right\| \leq \mathrm{M}_{\omega} \tag{17}
\end{equation*}
$$

always holds.

- To solve problem (6), we propose two types of gradient based Jacobi-type algorithms (Jacobi-GLU and Jacobi-GQU) on $\mathbf{S L}_{m}(\mathbb{C})$, which can be seen as non-orthogonal analogues of the Jacobi-G algorithm on orthogonal or unitary group [20, 24, 44].
- We establish the global convergence and weak convergence of Jacobi-GLU and JacobiGQU algorithms using the Lojasiewicz gradient inequality, under the assumption that the iterates are bounded, that is, there exists $\mathrm{M}_{\boldsymbol{X}}>0$ such that

$$
\begin{equation*}
\left\|\boldsymbol{X}_{k}\right\| \leq \mathrm{M}_{\boldsymbol{X}} \tag{18}
\end{equation*}
$$

always holds. To our knowledge, this is the first time that the theoretical convergence is established for the Jacobi-type algorithms on $\mathbf{S L}_{m}(\mathbb{C})$.

- These algorithms and convergence results are summarized in Table 1.

Table 1. A summary of the proposed algorithms

| Model | Proposed <br> algorithms | Position | Elementary <br> functions | Global <br> convergence | Weak <br> convergence |
| :---: | :--- | :--- | :--- | :--- | :--- |
|  | BCD-GLU | Algorithm 1\&Subalgorithm 1a | $\varphi, \psi, \rho$ | Theorem 7.1 | Theorem 7.2 |
|  | BCD-GQU | Algorithm 1\&Subalgorithm 1b | $h, \varphi, \rho$ | Theorem 7.3 | Theorem 7.4 |
| Second reformulation $(6)$ <br> on $\mathbf{S L}_{m}(\mathbb{C})$ | Jacobi-GLU | Algorithm 2 | Algorithm 1\&Subalgorithm 1c | $\varphi, \rho$ | $\varphi, \psi, \rho$ |

Remark 1.1. This paper is based on the complex matrices, complex Stiefel manifold $\mathbf{S t}(m, n, \mathbb{C})$ and complex special linear group $\mathbf{S L}_{m}(\mathbb{C})$. In fact, all the algorithms and convergence results described in this paper also apply to the real case.
1.7. Organization. The paper is organized as follows. In Section 2, we recall the basics of first order geometries on the Stiefel manifold $\mathbf{S t}(m, n, \mathbb{C})$ and special linear group $\mathbf{S L}_{m}(\mathbb{C})$, as well as the Lojasiewicz gradient inequality. In Section 3, we show the details of how to use line search descent method to update the first block variable in $\mathbf{S t}(m, n, \mathbb{C})$. In Section 4, using four kinds of elementary transformations, we construct three classes of elementary transformations to update the second block variable in $\mathbf{S L}_{m}(\mathbb{C})$, which induce three BCD-G algorithms (BCDGLU, BCD-GQU and BCD-GU) to solve problem (5). In Section 5 and Section 6, we present the details of these four kinds of elementary transformations. In Section 7, we prove our main results about the global and weak convergence of BCD-G algorithms. We also present two

[^4]gradient based Jacobi-type algorithms (Jacobi-GLU and Jacobi-GQU) to solve problem (6), and prove their global and weak convergence as well. In Section 8, some experiments are conducted to show the efficiency of the proposed algorithms. Section 9 concludes this paper with some final remarks and possible future work.

## 2. Geometries on $\mathbf{S t}(m, n, \mathbb{C})$ and $\mathbf{S L}_{m}(\mathbb{C})$

2.1. Notations. Let $1 \leq m \leq n$. For a complex matrix $\boldsymbol{X} \in \mathbb{C}^{n \times m}$ and a complex number $z \in \mathbb{C}$, we write the real and imaginary parts as $\boldsymbol{X}=\boldsymbol{X}^{\Re}+i \boldsymbol{X}^{\Im}$ and $z=\Re(z)+i \Im(z)$, respectively. For $\boldsymbol{X}, \boldsymbol{Y} \in \mathbb{C}^{n \times m}$, we introduce the following real-valued inner product

$$
\begin{equation*}
\langle\boldsymbol{X}, \boldsymbol{Y}\rangle \stackrel{\text { def }}{=}\left\langle\boldsymbol{X}^{\Re}, \boldsymbol{Y}^{\Re}\right\rangle+\left\langle\boldsymbol{X}^{\Im}, \boldsymbol{Y}^{\Im}\right\rangle=\Re\left(\operatorname{tr}\left(\boldsymbol{X}^{\mathrm{H}} \boldsymbol{Y}\right)\right) \tag{19}
\end{equation*}
$$

which makes $\mathbb{C}^{n \times m}$ a real Euclidean space of dimension $2 n m$. Let $h: \mathbb{C}^{n \times m} \rightarrow \mathbb{R}$ be a differentiable function and $\boldsymbol{X} \in \mathbb{C}^{n \times m}$. We denote by $\frac{\partial h}{\partial \boldsymbol{X}^{\Re}}, \frac{\partial h}{\partial \boldsymbol{X}^{\Im}} \in \mathbb{R}^{n \times m}$ the matrix Euclidean derivatives of $h$ with respect to real and imaginary parts of $\boldsymbol{X}$. The Wirtinger derivatives [1, 10, 22] are defined as

$$
\frac{\partial h}{\partial \boldsymbol{X}^{*}}:=\frac{1}{2}\left(\frac{\partial h}{\partial \boldsymbol{X}^{\Re}}+i \frac{\partial h}{\partial \boldsymbol{X}^{\Im}}\right), \quad \frac{\partial h}{\partial \boldsymbol{X}}:=\frac{1}{2}\left(\frac{\partial h}{\partial \boldsymbol{X}^{\Re}}-i \frac{\partial h}{\partial \boldsymbol{X}^{\Im}}\right)
$$

Then the Euclidean gradient of $h$ with respect to the inner product (19) becomes

$$
\begin{equation*}
\nabla h(\boldsymbol{X})=\frac{\partial h}{\partial \boldsymbol{X}^{\Re}}+i \frac{\partial h}{\partial \boldsymbol{X}^{\Im}}=2 \frac{\partial h}{\partial \boldsymbol{X}^{*}} \tag{20}
\end{equation*}
$$

For real matrices $\boldsymbol{X}, \boldsymbol{Y} \in \mathbb{R}^{n \times m}$, we see that (19) becomes the standard inner product, and (20) becomes the standard Euclidean gradient. We denote by $\mathbb{S}_{n-1} \subseteq \mathbb{R}^{n}$ the unit sphere.
2.2. Riemannian gradient on $\operatorname{St}(m, n, \mathbb{C})$. For $\boldsymbol{X} \in \mathbb{C}^{m \times m}$, we denote

$$
\operatorname{sym}(\boldsymbol{X})=\frac{1}{2}\left(\boldsymbol{X}+\boldsymbol{X}^{\mathrm{H}}\right), \quad \operatorname{skew}(\boldsymbol{X})=\frac{1}{2}\left(\boldsymbol{X}-\boldsymbol{X}^{\mathrm{H}}\right)
$$

Let $\mathbf{T}_{\boldsymbol{U}} \mathbf{S t}(m, n, \mathbb{C})$ be the tangent space to $\mathbf{S t}(m, n, \mathbb{C})$ at a point $\boldsymbol{U}$. Let $\boldsymbol{U}_{\perp} \in \mathbb{C}^{n \times(n-m)}$ be an orthogonal complement of $\boldsymbol{U}$, that is, $\left[\boldsymbol{U}, \boldsymbol{U}_{\perp}\right] \in \mathbb{C}^{n \times n}$ is a unitary matrix. By [33, Definition $6]$, we know that

$$
\mathbf{T}_{\boldsymbol{U}} \mathbf{S t}(m, n, \mathbb{C})=\left\{\boldsymbol{Z} \in \mathbb{C}^{n \times m}, \boldsymbol{Z}=\boldsymbol{U} \boldsymbol{A}+\boldsymbol{U}_{\perp} \boldsymbol{B}, \boldsymbol{A} \in \mathbb{C}^{m \times m}, \boldsymbol{A}^{\mathrm{H}}+\boldsymbol{A}=0, \boldsymbol{B} \in \mathbb{C}^{(n-m) \times m}\right\}
$$

which is a $\left(2 n m-m^{2}\right)$-dimensional vector space. The orthogonal projection of $\xi \in \mathbb{C}^{n \times m}$ to $\mathbf{T}_{\boldsymbol{U}} \mathbf{S t}(m, n, \mathbb{C})$ is

$$
\begin{equation*}
\operatorname{Proj}_{\boldsymbol{U}} \xi=\left(\boldsymbol{I}_{n}-\boldsymbol{U} \boldsymbol{U}^{\mathrm{H}}\right) \xi+\boldsymbol{U} \operatorname{skew}\left(\boldsymbol{U}^{\mathrm{H}} \xi\right)=\xi-\boldsymbol{U} \operatorname{sym}\left(\boldsymbol{U}^{\mathrm{H}} \xi\right) \tag{21}
\end{equation*}
$$

We denote $\operatorname{Proj}_{\boldsymbol{U}}^{\perp} \xi \stackrel{\text { def }}{=} \xi-\operatorname{Proj}_{\boldsymbol{U}} \xi$. Note that $\mathbf{S t}(m, n, \mathbb{C})$ is an embedded submanifold of the Euclidean space $\mathbb{C}^{n \times m}$. By (21), we have the Riemannian gradient of $h$ at $\boldsymbol{U}$ as:

$$
\begin{equation*}
\operatorname{grad} h(\boldsymbol{U})=\operatorname{Proj}_{\boldsymbol{U}} \nabla h(\boldsymbol{U})=\nabla h(\boldsymbol{U})-\boldsymbol{U} \operatorname{sym}\left(\boldsymbol{U}^{\mathrm{H}} \nabla h(\boldsymbol{U})\right) \tag{22}
\end{equation*}
$$

By [3, Example 5.4.2], the exponential map at $\boldsymbol{U}$ is defined as

$$
\begin{aligned}
\operatorname{Exp}_{\boldsymbol{U}}: \mathbf{T}_{\boldsymbol{U}} \operatorname{St}(m, n, \mathbb{C}) & \longrightarrow \operatorname{St}(m, n, \mathbb{C}) \\
\boldsymbol{Z} & \longmapsto[\boldsymbol{U}, \boldsymbol{Z}] \exp \left(\left[\begin{array}{cc}
\boldsymbol{U}^{H} \boldsymbol{Z} & -\boldsymbol{Z}^{H} \boldsymbol{Z} \\
\boldsymbol{I}_{m} & \boldsymbol{U}^{H} \boldsymbol{Z}
\end{array}\right]\right)\left[\begin{array}{c}
\exp \left(-\boldsymbol{U}^{H} \boldsymbol{Z}\right) \\
\mathbf{0}_{m \times m}
\end{array}\right] .
\end{aligned}
$$

2.3. Matrix groups. Let $\mathfrak{s l}_{m}(\mathbb{C}) \stackrel{\text { def }}{=}\left\{\boldsymbol{X} \in \mathbb{C}^{m \times m}, \operatorname{tr}(\boldsymbol{X})=0\right\}$. Then the tangent space to $\mathbf{S L}_{m}(\mathbb{C})$ at a point $\boldsymbol{X} \in \mathbf{S L}_{m}(\mathbb{C})$ can be constructed [7, Eq. (3.7),(3.8)] by $\mathbf{T}_{\boldsymbol{X}} \mathbf{S L}_{m}(\mathbb{C})=$ $\left\{\boldsymbol{X} \Omega, \Omega \in \mathfrak{s l}_{m}(\mathbb{C})\right\}$. Let $\mathbf{S U}_{m} \subseteq \mathbb{C}^{m \times m}$ be the special unitary group. Let

$$
\mathfrak{s u}_{m}(\mathbb{C}) \stackrel{\text { def }}{=}\left\{\boldsymbol{X} \in \mathbb{C}^{m \times m}, \boldsymbol{X}^{\boldsymbol{H}}=-\boldsymbol{X}, \operatorname{tr}(\boldsymbol{X})=0\right\}
$$

Then the tangent space to $\mathbf{S U}_{m}$ at a point $\boldsymbol{X} \in \mathbf{S U}_{m}$ can be constructed [7, Eq. (3.15)] by $\mathbf{T}_{\boldsymbol{X}} \mathbf{S U}_{m}=\left\{\boldsymbol{X} \Omega, \Omega \in \mathfrak{s l}_{m}(\mathbb{C})\right\}$.

A matrix $\boldsymbol{X} \in \mathbb{C}^{m \times m}$ is said to be upper triangular if $\mathrm{X}_{i j}=0$ for $i>j$. An upper triangular matrix $\boldsymbol{X}$ is said to be unipotent if it satisfies $\mathrm{X}_{i i}=1$ for $1 \leq i \leq m$. Let $\mathbf{U T}_{m}(\mathbb{C}) \subseteq \mathbf{G} \mathbf{L}_{m}(\mathbb{C})$ be the upper triangular subgroup, that is

$$
\mathbf{U T}_{m}(\mathbb{C}) \stackrel{\text { def }}{=}\left\{\boldsymbol{X} \in \mathbf{G L}_{m}(\mathbb{C}), \boldsymbol{X} \text { is upper triangular }\right\} .
$$

Let $\mathbf{S U T}_{m}(\mathbb{C}) \subseteq \mathbf{G L}_{m}(\mathbb{C})$ be the upper unipotent subgroup, that is

$$
\mathbf{S U T}_{m}(\mathbb{C}) \stackrel{\text { def }}{=}\left\{\boldsymbol{X} \in \mathbf{G L}_{m}(\mathbb{C}), \boldsymbol{X} \text { is upper triangular and unipotent }\right\} .
$$

A matrix $\boldsymbol{X} \in \mathbb{C}^{m \times m}$ is said to be strictly upper triangular if $\mathrm{X}_{i j}=0$ for $i \geq j$. Let $\mathfrak{s u t}_{m}(\mathbb{C}) \subseteq$ $\mathbb{C}^{m \times m}$ be the set of strictly upper triangular matrices. Then the tangent space to $\mathbf{S U T}_{m}(\mathbb{C})$ at a point $\boldsymbol{X} \in \mathbf{S U T}_{m}(\mathbb{C})$ can be constructed [7, Eq. (3.11)], [4, Section 6.4] by $\mathbf{T}_{\boldsymbol{X}} \mathbf{S U T}_{m}(\mathbb{C})=$ $\left\{\boldsymbol{X} \Omega, \Omega \in \mathfrak{s u t}_{m}(\mathbb{C})\right\}$. For the case $n=2$,

$$
\mathbf{S U T}_{2}(\mathbb{C})=\left\{\left[\begin{array}{ll}
1 & z \\
0 & 1
\end{array}\right], z \in \mathbb{C}\right\}, \quad \mathfrak{s u t}_{2}(\mathbb{C})=\left\{\left[\begin{array}{ll}
0 & z \\
0 & 0
\end{array}\right], z \in \mathbb{C}\right\} .
$$

Similary, we let $\mathbf{L T}_{m}(\mathbb{C}) \subseteq \mathbb{C}^{m \times m}$ be the lower triangular subgroup, $\mathbf{S L T}_{m}(\mathbb{C}) \subseteq \mathbb{C}^{m \times m}$ be the lower unipotent subgroup, and $\mathfrak{s l t}_{n}(\mathbb{C}) \subseteq \mathbb{C}^{m \times m}$ be the set of strictly lower triangular matrices. Then the tangent space to $\mathbf{S L T}_{m}(\mathbb{C})$ at a point $\boldsymbol{X} \in \mathbf{S L T}_{m}(\mathbb{C})$ can be constructed by $\mathbf{T}_{\boldsymbol{X}} \mathbf{S L T}_{m}(\mathbb{C})=\left\{\boldsymbol{X} \Omega, \Omega \in \mathfrak{s i t}_{m}(\mathbb{C})\right\}$.

A diagonal matrix $\boldsymbol{X} \in \mathbb{C}^{m \times m}$ is said to be a diagonal transformation if the product of all the diagonal elements is equal to 1 . Let $\mathbf{D}_{m}(\mathbb{C}) \subseteq \mathbf{G} \mathbf{L}_{m}(\mathbb{C})$ be the set of diagonal transformation matrices. Let $\mathfrak{d}_{m}(\mathbb{C}) \subseteq \mathbb{C}^{m \times m}$ be the set of diagonal traceless matrices. Then the tangent space to $\mathbf{D}_{m}(\mathbb{C})$ at a point $\boldsymbol{X} \in \mathbf{D}_{m}(\mathbb{C})$ can be constructed by $\mathbf{T}_{\boldsymbol{X}} \mathbf{D}_{m}(\mathbb{C})=\left\{\boldsymbol{X} \Omega, \Omega \in \mathfrak{d}_{m}(\mathbb{C})\right\}$. For the case $n=2$,

$$
\mathbf{D}_{2}(\mathbb{C})=\left\{\left[\begin{array}{cc}
z & 0 \\
0 & \frac{1}{z}
\end{array}\right], z \in \mathbb{C}_{*}\right\}, \quad \mathfrak{D}_{2}(\mathbb{C})=\left\{\left[\begin{array}{cc}
z & 0 \\
0 & -z
\end{array}\right], z \in \mathbb{C}\right\}
$$

2.4. Riemannian gradient on $\mathbf{S L}_{m}(\mathbb{C})$. For tangent vectors $\xi, \eta \in \mathbf{T}_{\boldsymbol{X}} \mathbf{S L}_{m}(\mathbb{C})$, we use the left invariant [3], [4, Eq. (6.2)] Riemannian metric

$$
\langle\xi, \eta\rangle_{\boldsymbol{X}} \stackrel{\text { def }}{=}\left\langle\boldsymbol{X}^{-1} \xi, \boldsymbol{X}^{-1} \eta\right\rangle=\Re\left(\operatorname{tr}\left(\xi^{\mathrm{H}}\left(\boldsymbol{X} \boldsymbol{X}^{\boldsymbol{H}}\right)^{-1} \eta\right)\right) .
$$

Let $g: \mathbf{S L}_{m}(\mathbb{C}) \longrightarrow \mathbb{R}^{+}$be a differentiable cost function. Then the Riemannian gradient of $g$ at $\boldsymbol{X} \in \mathbf{S L}_{m}(\mathbb{C})$ is the orthogonal projection [4, Lemma 6.2] of its Euclidian gradient $\nabla g(\boldsymbol{X})$ to $\mathbf{T}_{\boldsymbol{X}} \mathbf{S L}_{m}(\mathbb{C})$, that is,

$$
\begin{equation*}
\operatorname{grad} g(\boldsymbol{X})=\boldsymbol{X}\left(\boldsymbol{X}^{\mathrm{H}} \nabla g(\boldsymbol{X})-\frac{\operatorname{tr}\left(\boldsymbol{X}^{\mathrm{H}} \nabla g(\boldsymbol{X})\right)}{n} \boldsymbol{I}_{n}\right) \tag{23}
\end{equation*}
$$

We denote $\boldsymbol{\Lambda}(\boldsymbol{X}) \stackrel{\text { def }}{=} \boldsymbol{X}^{-1} \operatorname{grad} g(\boldsymbol{X}) \in \mathfrak{s l}_{m}(\mathbb{C})$ for $\boldsymbol{X} \in \mathbf{S L}_{m}(\mathbb{C})$, which will be frequently used in this paper.

In what follows, we will use the following exponential map

$$
\operatorname{Exp}_{\boldsymbol{X}}: \mathbf{T}_{\boldsymbol{X}} \mathbf{S L}_{m}(\mathbb{C}) \rightarrow \mathbf{S L}_{m}(\mathbb{C}), \quad \boldsymbol{X} \Omega \mapsto \boldsymbol{X} \exp \left(\Omega^{*}\right) \exp \left(\Omega-\Omega^{*}\right)
$$

where $\exp (\Omega)$ is the matrix exponential function $[3,7,19]$. For any $\Delta \in \mathbf{T}_{\boldsymbol{X}} \mathbf{S L}_{m}(\mathbb{C})$, we have the following relationship between $\operatorname{Exp}_{\boldsymbol{X}}$ and the Riemannian gradient

$$
\begin{equation*}
\langle\boldsymbol{\Delta}, \operatorname{grad} g(\boldsymbol{X})\rangle_{\boldsymbol{X}}=\left.\left(\frac{d}{d t} g\left(\operatorname{Exp}_{\boldsymbol{X}}(t \boldsymbol{\Delta})\right)\right)\right|_{t=0} \tag{24}
\end{equation*}
$$

### 2.5. Conditions for convergence analysis.

2.5.1. Łojasiewicz gradient inequality. In this subsection, we present some results about the Łojasiewicz gradient inequality $[23,30,2,43]$. These results were used in $[25,44]$ to prove the global convergence of Jacobi-G algorithms on the orthogonal and unitary groups.

Definition 2.1 ([40, Definition 2.1]). Let $\mathcal{M} \subseteq \mathbb{R}^{n}$ be a Riemannian submanifold, and $f$ : $\mathcal{M} \rightarrow \mathbb{R}$ be a differentiable function. The function $f: \mathcal{M} \rightarrow \mathbb{R}$ is said to satisfy a Eojasiewicz gradient inequality at $\boldsymbol{x} \in \mathcal{M}$, if there exist $\sigma>0, \zeta \in\left(0, \frac{1}{2}\right]$ and a neighborhood $\mathcal{U}$ in $\mathcal{M}$ of $\boldsymbol{x}$ such that for all $\boldsymbol{y} \in \mathcal{U}$, it follows that

$$
\begin{equation*}
|f(\boldsymbol{y})-f(\boldsymbol{x})|^{1-\zeta} \leq \sigma\|\operatorname{grad} f(\boldsymbol{y})\| . \tag{25}
\end{equation*}
$$

Lemma 2.2 ([40, Proposition 2.2]). Let $\mathcal{M} \subseteq \mathbb{R}^{n}$ be an analytic submanifold ${ }^{7}$ and $f: \mathcal{M} \rightarrow \mathbb{R}$ be a real analytic function. Then for any $\boldsymbol{x} \in \mathcal{M}, f$ satisfies a Lojasiewicz gradient inequality (25) in the $\delta$-neighborhood of $\boldsymbol{x}$, for some ${ }^{8} \delta, \sigma>0$ and $\zeta \in\left(0, \frac{1}{2}\right]$.

Theorem 2.3 ([40, Theorem 2.3]). Let $\mathcal{M} \subseteq \mathbb{R}^{n}$ be an analytic submanifold and $\left\{\boldsymbol{x}_{k}\right\}_{k \geq 1} \subseteq \mathcal{M}$. Suppose that $f$ is real analytic and, for large enough $k$,
(i) there exists $\sigma>0$ such that

$$
f\left(\boldsymbol{x}_{k+1}\right)-f\left(\boldsymbol{x}_{k}\right) \geq \sigma\left\|\operatorname{grad} f\left(\boldsymbol{x}_{k}\right)\right\|\left\|\boldsymbol{x}_{k+1}-\boldsymbol{x}_{k}\right\| ;
$$

(ii) $\operatorname{grad} f\left(\boldsymbol{x}_{k}\right)=0$ implies that $\boldsymbol{x}_{k+1}=\boldsymbol{x}_{k}$.

Then any accumulation point $\boldsymbol{x}_{*}$ of $\left\{\boldsymbol{x}_{k}\right\}_{k \geq 1}$ must be the only limit point.

[^5]2.5.2. Conditions for global convergence. The following result about the global convergence of Algorithm 1 is a direct consequence of Theorem 2.3 and condition (10).

Proposition 2.4. Suppose that the function $f$ in (8) is real analytic and the sequence $\left\{\boldsymbol{\omega}_{k}\right\}_{k \geq 1}$ produced by Algorithm 1 satisfies that, for large enough $k$,
(i) there exists $\sigma>0$ such that

$$
\begin{equation*}
f_{t_{k}}\left(\boldsymbol{\omega}_{k-1}\right)-f_{t_{k}}\left(\boldsymbol{\omega}_{k}\right) \geq \sigma\left\|\operatorname{grad} f_{t_{k}}\left(\boldsymbol{\omega}_{k-1}\right)\right\|\left\|\boldsymbol{\omega}_{k}-\boldsymbol{\omega}_{k-1}\right\| ; \tag{26}
\end{equation*}
$$

(ii) grad $f_{t_{k}}\left(\boldsymbol{\omega}_{k-1}\right)=0$ implies that $\boldsymbol{\omega}_{k}=\boldsymbol{\omega}_{k-1}$.

Then any accumulation point $\boldsymbol{\omega}_{*}$ of the sequence $\left\{\boldsymbol{\omega}_{k}\right\}_{k \geq 1}$ must be the only limit point.
2.5.3. Conditions for weak convergence. Since the special linear group $\mathbf{S L}_{m}(\mathbb{C})$ is not compact, the iterates $\left\{\boldsymbol{\omega}_{k}\right\}_{k \geq 1}$ produced by Algorithm 1 may have no accumulation point. However, if there exists an accumulation point, we have the following result about its weak convergence, which can be proved easily by condition (10) and the fact that $f(\omega) \geq 0$.

Lemma 2.5. In Algorithm 1 for cost function (8), if there exists $\eta>0$ such that

$$
\begin{equation*}
f\left(\boldsymbol{\omega}_{k-1}\right)-f\left(\boldsymbol{\omega}_{k}\right) \geq \eta\left\|\operatorname{grad} f_{t_{k}}\left(\boldsymbol{\omega}_{k-1}\right)\right\|^{2} \tag{27}
\end{equation*}
$$

always holds, then $\lim _{k \rightarrow \infty} \operatorname{grad} f\left(\boldsymbol{\omega}_{k-1}\right)=0$. In particular, if $\boldsymbol{\omega}_{*}$ is an accumulation point of the iterates $\left\{\boldsymbol{\omega}_{k}\right\}_{k \geq 1}$, then $\boldsymbol{\omega}_{*}$ is a stationary point of $f$.

## 3. Line search descent method on $\mathbf{S t}(m, n, \mathbb{C})$

Let $f$ be the cost function (5). Let $\boldsymbol{\omega}_{k-1}=\left(\boldsymbol{U}_{k-1}, \boldsymbol{X}_{k-1}\right)$ and $p=f_{1, \boldsymbol{X}_{k-1}}$ be the first restricted function. Denote $\boldsymbol{X}=\boldsymbol{X}_{k-1}$ for simplicity. Then $p$ can be expressed as

$$
\begin{equation*}
p: \mathbf{S t}(m, n, \mathbb{C}) \rightarrow \mathbb{R}^{+}, \quad \boldsymbol{U} \mapsto \sum_{\ell=1}^{L} \mu_{\ell}\left\|\operatorname{offdiag}\left\{\boldsymbol{W}^{(\ell)}\right\}\right\|^{2}, \tag{28}
\end{equation*}
$$

where $\boldsymbol{W}^{(\ell)}=\boldsymbol{X}^{\star} U^{\star} \boldsymbol{A}^{(\ell)} \boldsymbol{U} \boldsymbol{X}$.
3.1. Riemannian gradient. We first present a lemma, which can be obtained by direct calculations. This result will help us to obtain the Riemannian gradient of $p$ in (28).

Lemma 3.1. Let $\boldsymbol{A} \in \mathbb{C}^{n \times n}$ and the function $\tilde{p}$ be defined as

$$
\tilde{p}: \mathbb{C}^{n \times m} \rightarrow \mathbb{R}^{+}, \quad \boldsymbol{Y} \mapsto \| \text { offdiag }\{\boldsymbol{W}\} \|^{2},
$$

where $\boldsymbol{W}=\boldsymbol{Y}^{\star} \boldsymbol{A} \boldsymbol{Y}$. Let $\boldsymbol{V}=\boldsymbol{A} \boldsymbol{Y}=\left[\boldsymbol{v}_{1}, \cdots, \boldsymbol{v}_{m}\right]$ and $\boldsymbol{V}^{\prime}=\boldsymbol{A}^{\boldsymbol{\wedge}} \boldsymbol{Y}=\left[\boldsymbol{v}_{1}^{\prime}, \cdots, \boldsymbol{v}_{m}^{\prime}\right]$. Denote $\boldsymbol{Y}=\left[\boldsymbol{y}_{1}, \cdots, \boldsymbol{y}_{m}\right]$. Then the Euclidean gradient is

$$
\nabla \tilde{p}(\boldsymbol{Y})=2\left(\sum_{j \neq 1} \boldsymbol{v}_{j} \boldsymbol{v}_{j}^{\mathrm{H}} \boldsymbol{y}_{1}, \cdots, \sum_{j \neq m} \boldsymbol{v}_{j} \boldsymbol{v}_{j}^{\mathrm{H}} \boldsymbol{y}_{m}\right)+2\left(\sum_{j \neq 1} \boldsymbol{v}_{j}^{\prime}\left(\boldsymbol{v}_{j}^{\prime}\right)^{\mathrm{H}} \boldsymbol{y}_{1}, \cdots, \sum_{j \neq m} \boldsymbol{v}_{j}^{\prime}\left(\boldsymbol{v}_{j}^{\prime}\right)^{\mathrm{H}} \boldsymbol{y}_{m}\right) .
$$

In particular, it satisfies

$$
\begin{equation*}
\boldsymbol{Y} \boldsymbol{Y}^{\mathrm{H}} \nabla \tilde{p}(\boldsymbol{Y})=2 \boldsymbol{Y} \Upsilon(\boldsymbol{W}), \tag{29}
\end{equation*}
$$

where

$$
\Upsilon(\boldsymbol{W}) \stackrel{\text { def }}{=}\left\{\begin{array}{cl}
\boldsymbol{W} \operatorname{offdiag}\{\boldsymbol{W}\}^{\mathrm{H}}+\boldsymbol{W}^{H} \operatorname{offdiag}\{\boldsymbol{W}\}, & \text { if }(\cdot)^{\star}=(\cdot)^{\mathrm{H}} ; \\
\boldsymbol{W}^{*} \operatorname{offdiag}\{\boldsymbol{W}\}^{\top}+\boldsymbol{W}^{H} \operatorname{offdiag}\{\boldsymbol{W}\}, & \text { if }(\cdot)^{\star}=(\cdot)^{\mathrm{T}} .
\end{array}\right.
$$

Lemma 3.2. Let $\boldsymbol{W}^{(\ell)}$ and the function $p$ be as in (28). Then the Euclidean gradient satisfies

$$
\begin{equation*}
\boldsymbol{U}^{\mathrm{H}} \nabla p(\boldsymbol{U})=2\left(\boldsymbol{X}^{\mathrm{H}}\right)^{-1} \sum_{\ell=1}^{L} \mu_{\ell} \Upsilon\left(\boldsymbol{W}^{(\ell)}\right) \boldsymbol{X}^{\mathrm{H}} . \tag{30}
\end{equation*}
$$

Proof. By the product rule, we see that $\nabla p(\boldsymbol{U})=\nabla \tilde{p}(\boldsymbol{U} \boldsymbol{X}) \boldsymbol{X}^{\mathrm{H}}$. Then, by (29), we have

$$
\boldsymbol{X} \boldsymbol{X}^{\mathrm{H}} \boldsymbol{U}^{\mathrm{H}} \nabla p(\boldsymbol{U})=\boldsymbol{U}^{\mathrm{H}} \boldsymbol{U} \boldsymbol{X}(\boldsymbol{U} \boldsymbol{X})^{\mathrm{H}} \nabla \tilde{p}(\boldsymbol{U} \boldsymbol{X}) \boldsymbol{X}^{\mathrm{H}}=2 \boldsymbol{U}^{\mathrm{H}} \boldsymbol{U} \boldsymbol{X} \sum_{\ell=1}^{L} \mu_{\ell} \Upsilon\left(\boldsymbol{W}^{(\ell)}\right) \boldsymbol{X}^{\mathrm{H}} .
$$

Note that $\boldsymbol{X}$ is invertible. The proof is complete.
Now, by (30) and (22), we see that the Riemannian gradient of $p$ in (28) satisfies

$$
\boldsymbol{U}^{\mathrm{H}} \operatorname{grad} p(\boldsymbol{U})=\left(\boldsymbol{X}^{\mathrm{H}}\right)^{-1} \sum_{\ell=1}^{L} \mu_{\ell} \Upsilon\left(\boldsymbol{W}^{(\ell)}\right) \boldsymbol{X}^{\mathrm{H}}-\boldsymbol{X} \sum_{\ell=1}^{L} \mu_{\ell} \Upsilon\left(\boldsymbol{W}^{(\ell)}\right)^{\mathrm{H}} \boldsymbol{X}^{-1} .
$$

3.2. Line search descent method. In this paper, we adopt the line search descent method [2, 3, 37, 38, 39] on $\mathbf{S t}(m, n, \mathbb{C})$ to find $\boldsymbol{U}_{k}$ for the cost function (28). More precisely, we set

$$
\begin{equation*}
\boldsymbol{U}_{k}=\operatorname{Exp}_{\boldsymbol{U}_{k-1}}\left(t_{k-1} \boldsymbol{Z}_{k-1}\right), \tag{31}
\end{equation*}
$$

where $\boldsymbol{Z}_{k-1}$ is the search direction and $t_{k-1}$ is the step size. We always choose $\boldsymbol{Z}_{k-1}$ such that

$$
\begin{equation*}
\left\langle\operatorname{grad} p\left(\boldsymbol{U}_{k-1}\right), \boldsymbol{Z}_{k-1}\right\rangle_{\boldsymbol{U}_{k-1}} \leq-\delta_{s}\left\|\operatorname{grad} p\left(\boldsymbol{U}_{k-1}\right)\right\|\left\|\boldsymbol{Z}_{k-1}\right\|, \tag{32}
\end{equation*}
$$

with $0<\delta_{s}<1$. We say that $t_{k-1}$ satisfies the Armijo condition ${ }^{9}$, if

$$
\begin{equation*}
p\left(\boldsymbol{U}_{k}\right) \leq p\left(\boldsymbol{U}_{k-1}\right)+\delta_{w} t_{k-1}\left\langle\operatorname{grad} p\left(\boldsymbol{U}_{k-1}\right), \boldsymbol{Z}_{k-1}\right\rangle_{\boldsymbol{U}_{k-1}}, \tag{33}
\end{equation*}
$$

with $0<\delta_{w}<1$. We say that $t_{k-1}$ satisfies the curvature condition, if

$$
\begin{equation*}
\left\langle\operatorname{grad} p\left(\boldsymbol{U}_{k}\right), \mathbf{D E x p}_{\boldsymbol{U}_{k-1}}\left(t_{k-1} \boldsymbol{Z}_{k-1}\right)\left[\boldsymbol{Z}_{k-1}\right]\right\rangle_{\boldsymbol{U}_{k}} \geq c_{2}\left\langle\operatorname{grad} p\left(\boldsymbol{U}_{k-1}\right), \boldsymbol{Z}_{k-1}\right\rangle_{\boldsymbol{U}_{k-1}}, \tag{34}
\end{equation*}
$$

with $\delta_{w}<c_{2}<1$. The conditions (33) and (34) are known collectively as the Wolfe conditions. As in the Euclidean space case [37, Lemma 3.1], it was shown [38, 39] that we can always choose $t_{k-1}$ such that the conditions (33) and (34) are both satisfied. It is not difficult to see that there exists $\mathrm{M}_{e}>0$ such that

$$
\left\|\operatorname{Exp}_{\boldsymbol{U}}\left(\boldsymbol{Z}_{1}\right)-\operatorname{Exp}_{\boldsymbol{U}}\left(\boldsymbol{Z}_{2}\right)\right\| \leq \mathrm{M}_{e}\left\|\boldsymbol{Z}_{1}-\boldsymbol{Z}_{2}\right\|,
$$

for any $\boldsymbol{U} \in \mathbf{S t}(m, n, \mathbb{C})$ and $\boldsymbol{Z}_{1}, \boldsymbol{Z}_{2} \in \mathbf{T}_{\boldsymbol{U}} \boldsymbol{S t}(m, n, \mathbb{C})$. Then the next result follows directly.
Lemma 3.3. If we set $\boldsymbol{U}_{k}$ as in (31) such that the conditions (32) and (33) are both satisfied, then we have

$$
p\left(\boldsymbol{U}_{k-1}\right)-p\left(\boldsymbol{U}_{k}\right) \geq \delta_{s} \delta_{w}\left\|\operatorname{grad} p\left(\boldsymbol{U}_{k-1}\right)\right\|\left\|t_{k-1} \boldsymbol{Z}_{k-1}\right\| \geq \sigma_{p}\left\|\operatorname{grad} p\left(\boldsymbol{U}_{k-1}\right)\right\|\left\|\boldsymbol{U}_{k}-\boldsymbol{U}_{k-1}\right\|,
$$ where $\sigma_{p}=\frac{\delta_{s} \delta_{w}}{M_{e}}$.

[^6]Moreover, we have the next result, which is a simple corollary of the proof in [38, Theorem $2]$.

Lemma 3.4. If we set $\boldsymbol{U}_{k}$ as in (31) such that the conditions (32), (33) and (34) are all satisfied, then we have

$$
\begin{equation*}
p\left(\boldsymbol{U}_{k-1}\right)-p\left(\boldsymbol{U}_{k}\right) \geq \eta_{p}\left\|\operatorname{grad} p\left(\boldsymbol{U}_{k-1}\right)\right\|^{2} \tag{35}
\end{equation*}
$$

where $\eta_{p}>0$ is fixed.

## 4. GLU, GQU and GU classes of transformations on $\mathbf{S L}_{m}(\mathbb{C})$

4.1. Elementary functions. Let $(i, j)$ be a pair of indices with $1 \leq i<j \leq n$. Define a projection operator $\mathcal{P}_{i, j}: \mathbb{C}^{m \times m} \rightarrow \mathbb{C}^{2 \times 2}$ which extracts a submatrix of $\boldsymbol{X} \in \mathbb{C}^{m \times m}$ as follows:

$$
\mathcal{P}_{i, j}(\boldsymbol{X})=\left[\begin{array}{ll}
\mathrm{X}_{i i} & \mathrm{X}_{i j}  \tag{36}\\
\mathrm{X}_{j i} & \mathrm{X}_{j j}
\end{array}\right] .
$$

Conversely, we introduce an operator

$$
\mathcal{E}_{i, j}: \mathbb{C}^{2 \times 2} \rightarrow \mathbb{C}^{m \times m}, \Psi \mapsto \mathcal{E}_{i, j}(\boldsymbol{\Psi})
$$

where $\mathcal{E}_{i, j}(\boldsymbol{\Psi})$ is the identity matrix $\boldsymbol{I}_{m}$ except that $\mathcal{P}_{i, j}\left(\mathcal{E}_{i, j}(\boldsymbol{\Psi})\right)=\boldsymbol{\Psi}$. Now we define the following four elementary transformations in $\mathbf{S L}_{m}(\mathbb{C})$.

- $\boldsymbol{G}^{(i, j, \boldsymbol{\Psi})} \stackrel{\text { def }}{=} \mathcal{E}_{i, j}(\boldsymbol{\Psi})$ : Givens plane transformation for a matrix $\boldsymbol{\Psi} \in \mathbf{S U}_{2}$.
- $\boldsymbol{U}^{(i, j, \Psi)} \stackrel{\text { def }}{=} \mathcal{E}_{i, j}(\boldsymbol{\Psi})$ : plane upper triangular transformation for a matrix $\boldsymbol{\Psi} \in \mathbf{S U T}_{2}(\mathbb{C})$.
- $\boldsymbol{L}^{(i, j, \boldsymbol{\Psi})} \stackrel{\text { def }}{=} \mathcal{E}_{i, j}(\boldsymbol{\Psi})$ : plane lower triangular transformation for a matrix $\boldsymbol{\Psi} \in \mathbf{S L T}_{2}(\mathbb{C})$.
- $\boldsymbol{D}^{(i, j, \Psi)} \stackrel{\text { def }}{=} \mathcal{E}_{i, j}(\boldsymbol{\Psi})$ : plane diagonal transformation for a matrix $\boldsymbol{\Psi} \in \mathbf{D}_{2}(\mathbb{C})$.

Remark 4.1. These four elementary transformations have all been used in the literature. The Givens plane transformation $\boldsymbol{G}^{(i, j, \boldsymbol{\Psi})}$ was used very often in Jacobi-type algorithms for joint approximate diagonalization of matrices or tensors by orthogonal or non-orthogonal transformations [17, 24, 44, 4, 5, 41]. Plane triangular transformations $\boldsymbol{U}^{(i, j, \boldsymbol{\Psi})}$ and $\boldsymbol{L}^{(i, j, \boldsymbol{\Psi})}$ also appeared many times in the Jacobi-type algorithms on special linear group $\mathbf{S L}_{m}(\mathbb{C})$ or $\mathbf{S L}_{m}(\mathbb{R})$ $[4,5,34,35,36]$. In the real case, the diagonal transform $\boldsymbol{D}^{(i, j, \boldsymbol{\Psi})}$ was used in [42]. In this paper, using these four elementary transformations, we will construct three classes (GLU, GQU and GU ) in Sections 4.3 and 4.4, such that a gradient inequality (38) is always satisfied to establish the global convergence of Algorithm 1 in Section 7.

In this paper, as in [17, 44], we parameterize $\boldsymbol{\Psi} \in \mathbf{S U}_{2}$ as

$$
\boldsymbol{\Psi}=\boldsymbol{\Psi}\left(c, s_{1}, s_{2}\right)=\left[\begin{array}{cc}
c & -s \\
s^{*} & c
\end{array}\right]=\left[\begin{array}{cc}
c & -\left(s_{1}+i s_{2}\right) \\
s_{1}-i s_{2} & c
\end{array}\right]=\left[\begin{array}{cc}
\cos \theta & -\sin \theta \mathrm{e}^{i \phi} \\
\sin \theta \mathrm{e}^{-i \phi} & \cos \theta
\end{array}\right]
$$

where $\left(c, s_{1}, s_{2}\right) \in \mathbb{S}_{2}$ and $(\theta, \phi) \in \mathbb{R}^{2}$. Let $g: \mathbf{S L}_{m}(\mathbb{C}) \rightarrow \mathbb{R}^{+}$be a smooth function. Let $\boldsymbol{X} \in \mathbf{S L}_{m}(\mathbb{C})$ and $z=x+i y$. We define the following elementary functions:

$$
\begin{aligned}
h_{(i, j), \boldsymbol{X}}\left(c, s_{1}, s_{2}\right)= & h_{(i, j), \boldsymbol{X}}(\theta, \phi)=h_{(i, j), \boldsymbol{X}}(\boldsymbol{\Psi}) \stackrel{\text { def }}{=} g\left(\boldsymbol{X} \boldsymbol{G}^{(i, j, \boldsymbol{\Psi})}\right), \quad\left(c, s_{1}, s_{2}\right) \in \mathbb{S}_{2} ; \\
& \varphi_{(i, j), \boldsymbol{X}}(x, y)=\varphi_{(i, j), \boldsymbol{X}}(\boldsymbol{\Psi}) \stackrel{\text { def }}{=} g\left(\boldsymbol{X} \boldsymbol{U}^{(i, j, \boldsymbol{\Psi})}\right), \quad z \in \mathbb{C} ; \\
& \psi_{(i, j), \boldsymbol{X}}(x, y)=\psi_{(i, j), \boldsymbol{X}}(\boldsymbol{\Psi}) \stackrel{\text { def }}{=} g\left(\boldsymbol{X} \boldsymbol{L}^{(i, j, \boldsymbol{\Psi})}\right), \quad z \in \mathbb{C} ; \\
& \rho_{(i, j), \boldsymbol{X}}(x, y)=\rho_{(i, j), \boldsymbol{X}}(\boldsymbol{\Psi}) \stackrel{\text { def }}{=} g\left(\boldsymbol{X} \boldsymbol{D}^{(i, j, \boldsymbol{\Psi})}\right), \quad z \in \mathbb{C}_{*} .
\end{aligned}
$$

4.2. Derivatives of elementary functions. Denote $\boldsymbol{\Lambda}=\boldsymbol{\Lambda}(\boldsymbol{X})$, which is defined as in Section 2.4. Let $\mathcal{P}_{i, j}^{\top}: \mathbb{C}^{2 \times 2} \rightarrow \mathbb{C}^{m \times m}$ be the adjoint operator of projection operator $\mathcal{P}_{i, j}$ defined in (36), i.e.,

$$
\mathcal{P}_{i, j}^{\top}(\mathbf{\Psi})=\begin{gathered}
\\
i \\
i
\end{gathered}\left[\begin{array}{ccc}
\mathbf{0} & \vdots & \vdots \\
\cdots & \mathbf{0} \\
\cdots & \Psi_{i i} & \Psi_{i j} \\
\\
\cdots & \Psi_{j i} & \Psi_{j j} \\
\mathbf{0} & & \\
\mathbf{0}
\end{array}\right]
$$

for $\boldsymbol{\Psi} \in \mathbb{C}^{2 \times 2}$.
Lemma 4.2. The Riemannian gradients of the elementary functions defined in Section 4.1 at the identity matrix $\boldsymbol{I}_{2}$ can be expressed as follows:

- $\operatorname{grad} h_{(i, j), \boldsymbol{X}}\left(\boldsymbol{I}_{2}\right)=\left[\begin{array}{cc}\frac{i}{2} \Im\left(\boldsymbol{\Lambda}_{i i}-\boldsymbol{\Lambda}_{j j}\right) \\ -\frac{1}{2} \Re\left(\boldsymbol{\Lambda}_{i j}-\boldsymbol{\Lambda}_{j i}\right)+\frac{i}{2} \Im\left(\boldsymbol{\Lambda}_{i j}+\boldsymbol{\Lambda}_{j i}\right) & \left.\begin{array}{c}\frac{1}{2} \Re\left(\boldsymbol{\Lambda}_{i j}-\boldsymbol{\Lambda}_{j i}\right)+\frac{i}{2} \Im\left(\boldsymbol{\Lambda}_{i j}+\boldsymbol{\Lambda}_{j i}\right) \\ -\frac{i}{2} \Im\left(\boldsymbol{\Lambda}_{i i}-\boldsymbol{\Lambda}_{j j}\right)\end{array}\right] ; ; ~ ; ~\end{array}\right.$
- $\operatorname{grad} \varphi_{(i, j), \boldsymbol{X}}\left(\boldsymbol{I}_{2}\right)=\left[\begin{array}{cc}0 & \boldsymbol{\Lambda}_{i j} \\ 0 & 0\end{array}\right] ;$
- $\operatorname{grad} \psi_{(i, j), \boldsymbol{X}}\left(\boldsymbol{I}_{2}\right)=\left[\begin{array}{cc}0 & 0 \\ \boldsymbol{\Lambda}_{j i} & 0\end{array}\right]$;
- $\operatorname{grad} \rho_{(i, j), \boldsymbol{X}}\left(\boldsymbol{I}_{2}\right)=\left[\begin{array}{cc}\Re\left(\boldsymbol{\Lambda}_{i i}-\boldsymbol{\Lambda}_{j j}\right) & 0 \\ 0 & \Im\left(\boldsymbol{\Lambda}_{i i}-\boldsymbol{\Lambda}_{j j}\right)\end{array}\right]$.

Proof. We only prove the case of $h_{(i, j), \boldsymbol{X}}$. Other cases are similar. For any $\boldsymbol{\Delta} \in \mathfrak{s u}_{2}(\mathbb{C})=$ $\mathbf{T}_{\mathbf{I}_{2}} \mathbf{S U}_{2}$, by equation (24), we have

$$
\begin{aligned}
& \left\langle\boldsymbol{\Delta}, \operatorname{grad} h_{(i, j), \boldsymbol{X}}\left(\boldsymbol{I}_{2}\right)\right\rangle_{\boldsymbol{I}_{2}}=\left.\left(\frac{d}{d t} h_{(i, j), \boldsymbol{X}}\left(\operatorname{Exp}_{\boldsymbol{I}_{2}}(t \boldsymbol{\Delta})\right)\right)\right|_{t=0}=\left.\left(\frac{d}{d t} f\left(\boldsymbol{X} \boldsymbol{G}^{\left(i, j, \operatorname{Exp}_{\boldsymbol{I}_{2}}(\boldsymbol{\Delta} t)\right)}\right)\right)\right|_{t=0} \\
& =\left.\left(\frac{d}{d t} f\left(\operatorname{Exp}_{\boldsymbol{X}}\left(\boldsymbol{X}_{\mathcal{P}_{i, j}}^{\top}(\boldsymbol{\Delta}) t\right)\right)\right)\right|_{t=0}=\left\langle\boldsymbol{X} \mathcal{P}_{i, j}^{\top}(\boldsymbol{\Delta}), \operatorname{grad} f(\boldsymbol{X})\right\rangle_{\boldsymbol{X}}=\left\langle\boldsymbol{\Delta}, \mathcal{P}_{i, j}(\boldsymbol{\Lambda})\right\rangle_{\boldsymbol{I}_{2}} .
\end{aligned}
$$

The proof is complete.
The following lemma can be easily obtained by Lemma 4.2.

Lemma 4.3. The partial derivatives of the elementary functions defined in Section 4.1 satisfy

- $\partial h_{(i, j), \boldsymbol{X}}\left(\boldsymbol{I}_{2}\right) \stackrel{\text { def }}{=} \partial h_{(i, j), \boldsymbol{X}}(1,0,0)=\left[0,-\Re\left(\boldsymbol{\Lambda}_{i j}-\boldsymbol{\Lambda}_{j i}\right),-\Im\left(\boldsymbol{\Lambda}_{i j}+\boldsymbol{\Lambda}_{j i}\right)\right]^{\top} ;$
- $\partial \varphi_{(i, j), \boldsymbol{X}}\left(\boldsymbol{I}_{2}\right) \stackrel{\text { def }}{=} \partial \varphi_{(i, j), \boldsymbol{X}}(0,0)=\left[\Re\left(\boldsymbol{\Lambda}_{i j}\right), \Im\left(\boldsymbol{\Lambda}_{i j}\right)\right]^{\top}$;
- $\partial \psi_{(i, j), \boldsymbol{X}}\left(\boldsymbol{I}_{2}\right) \stackrel{\text { def }}{=} \partial \psi_{(i, j), \boldsymbol{X}}(0,0)=\left[\Re\left(\boldsymbol{\Lambda}_{j i}\right), \Im\left(\boldsymbol{\Lambda}_{j i}\right)\right]^{\top}$;
- $\partial \rho_{(i, j), \boldsymbol{X}}\left(\boldsymbol{I}_{2}\right) \stackrel{\text { def }}{=} \partial \rho_{(i, j), \boldsymbol{X}}(0,0)=\left[\Re\left(\boldsymbol{\Lambda}_{i i}-\boldsymbol{\Lambda}_{j j}\right), \Im\left(\boldsymbol{\Lambda}_{i i}-\boldsymbol{\Lambda}_{j j}\right)\right]^{\top}$.
4.3. GLU and GQU classes of transformations. Let $\boldsymbol{\omega}_{k-1}=\left(\boldsymbol{U}_{k-1}, \boldsymbol{X}_{k-1}\right)$ be the $(k-1)$ th iterate produced by Algorithm 1 , and $g: \mathbf{S L}_{m}(\mathbb{C}) \rightarrow \mathbb{R}$ be the restricted function $f_{2, \boldsymbol{U}_{k-1}}$ defined as in (9). Let $\left(i_{k}, j_{k}\right)$ be a pair of indices with $1 \leq i_{k}<j_{k} \leq m$. We denote that

$$
\begin{equation*}
h_{k}=h_{\left(i_{k}, j_{k}\right), \boldsymbol{X}_{k-1}}, \quad \varphi_{k}=\varphi_{\left(i_{k}, j_{k}\right), \boldsymbol{X}_{k-1}}, \quad \psi_{k}=\psi_{\left(i_{k}, j_{k}\right), \boldsymbol{X}_{k-1}}, \quad \rho_{k}=\rho_{\left(i_{k}, j_{k}\right), \boldsymbol{X}_{k-1}} \tag{37}
\end{equation*}
$$

Inspired by the equations (11), (12) and a similar idea as in [20, 24, 44], we now propose two classes of elementary transformations on $\mathbf{S L}_{m}(\mathbb{C})$. We call them the $G L U$ (based on LU decomposition) and $G Q U$ (based on QU decomposition) transformations ${ }^{10}$, respectively. Based on these two classes of elementary transformations, the subalgorithms to update $\boldsymbol{X}_{k}$ in Algorithm 1 are summarized in Subalgorithm 1a and Subalgorithm 1b, respectively.

Subalgorithm 1a: The subalgorithm to update $\boldsymbol{X}_{k}$ based on GLU transformations
Input: A fixed positive constant $0<\varepsilon<\sqrt{\frac{2}{3 m(m-1)}}$.
Output: New iterate $\boldsymbol{X}_{k}$.
Choose the index pair $\left(i_{k}, j_{k}\right)$ and $\nu_{k}$ such that

$$
\begin{equation*}
\left\|\partial \nu_{k}\left(\boldsymbol{I}_{2}\right)\right\| \geq \varepsilon\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\| \tag{38}
\end{equation*}
$$

where $\nu_{k}=\varphi_{k}, \psi_{k}$ or $\rho_{k}$;
Compute $\boldsymbol{\Psi}_{k}^{*}$ that minimizes the function $\nu_{k}$; Update $\boldsymbol{X}_{k}=\boldsymbol{X}_{k-1} \boldsymbol{V}_{k}$, where $\boldsymbol{V}_{k}=\boldsymbol{U}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}, \boldsymbol{L}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}$ or $\boldsymbol{D}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}$.

Subalgorithm 1b: The subalgorithm to update $\boldsymbol{X}_{k}$ based on GQU transformations
Input: A fixed positive constant $0<\varepsilon<\sqrt{\frac{3+\sqrt{5}}{3 m(m-1)}}$.
Output: New iterate $\boldsymbol{X}_{k}$. Choose the index pair $\left(i_{k}, j_{k}\right)$ and $\nu_{k}$ satisfying (38), where $\nu_{k}=h_{k}, \varphi_{k}$ or $\rho_{k}$; Compute $\boldsymbol{\Psi}_{k}^{*}$ that minimizes the function $\nu_{k}$; Update $\boldsymbol{X}_{k}=\boldsymbol{X}_{k-1} \boldsymbol{V}_{k}$, where $\boldsymbol{V}_{k}=\boldsymbol{G}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}, \boldsymbol{U}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}$ or $\boldsymbol{D}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}$.

Proposition 4.4. In Subalgorithm 1a and Subalgorithm 1b, we can always choose an index pair $\left(i_{k}, j_{k}\right)$ and an elementary function $\nu_{k}(x)$ such that the inequality (38) is satisfied.

We need a simple lemma before the proof of Proposition 4.4.

[^7]Lemma 4.5. (i) Let $z_{1}, z_{2} \in \mathbb{C}$. Then

$$
\left|z_{1}-z_{2}\right|^{2}+\left|z_{2}\right|^{2} \geq \frac{3-\sqrt{5}}{2}\left(\left|z_{1}\right|^{2}+\left|z_{2}\right|^{2}\right)
$$

(ii) Let $\left\{z_{i}\right\}_{1 \leq i \leq m} \subseteq \mathbb{C}$ satisfy that $\sum_{1 \leq i \leq m} z_{i}=0$. Then

$$
\sum_{1 \leq i<j \leq m}\left|z_{i}-z_{j}\right|^{2}=m \sum_{1 \leq i \leq m}\left|z_{i}\right|^{2}
$$

Proof of Proposition 4.4. We first prove the existence of such an index pair $\left(i_{k}, j_{k}\right)$ and elementary function $\nu_{k}$ in Subalgorithm 1a. By Lemma 4.3, Lemma 4.5 and $\boldsymbol{\Lambda}=\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right) \in \mathfrak{s l}_{m}(\mathbb{C})$, we get that

$$
\begin{aligned}
& \sum_{1 \leq i_{k}<j_{k} \leq m}\left(\left\|\partial \varphi_{k}\left(\boldsymbol{I}_{2}\right)\right\|^{2}+\left\|\partial \psi_{k}\left(\boldsymbol{I}_{2}\right)\right\|^{2}+\left\|\partial \rho_{k}\left(\boldsymbol{I}_{2}\right)\right\|^{2}\right) \\
& =\sum_{1 \leq i_{k}<j_{k} \leq m}\left(\left|\Lambda_{i_{k} j_{k}}\right|^{2}+\left|\Lambda_{j_{k} i_{k}}\right|^{2}\right)+m \sum_{1 \leq i_{k} \leq m}\left|\Lambda_{i_{k} i_{k}}\right|^{2} \geq\|\boldsymbol{\Lambda}\|^{2} .
\end{aligned}
$$

Therefore, there exist an index pair $\left(i_{k}, j_{k}\right)$ and $\nu_{k}(x)$ such that

$$
\left.\frac{3}{2} m(m-1) \right\rvert\,\left\|\partial \psi_{k}\left(\boldsymbol{I}_{2}\right)\right\|^{2} \geq\|\boldsymbol{\Lambda}\|^{2}
$$

Next, we prove the existence in Subalgorithm 1b. Similarly, we get that

$$
\begin{aligned}
& \sum_{1 \leq i_{k}<j_{k} \leq m}\left(\left\|\partial h_{k}\left(\boldsymbol{I}_{2}\right)\right\|^{2}+\left\|\partial \varphi_{k}\left(\boldsymbol{I}_{2}\right)\right\|^{2}+\left\|\partial \rho_{k}\left(\boldsymbol{I}_{2}\right)\right\|^{2}\right) \\
= & \sum_{1 \leq i_{k}<j_{k} \leq m}\left(\left|\Lambda_{i_{k} j_{k}}^{*}-\Lambda_{j_{k} i_{k}}\right|^{2}+\left|\Lambda_{i_{k} j_{k}}\right|^{2}+\left|\Lambda_{i_{k} i_{k}}-\Lambda_{j_{k} j_{k}}\right|^{2}\right) \\
\geq & \frac{3-\sqrt{5}}{2} \sum_{1 \leq i_{k}<j_{k} \leq m}\left(\left|\Lambda_{i_{k} j_{k}}\right|^{2}+\left|\Lambda_{j_{k} i_{k}}\right|^{2}\right)+m \sum_{1 \leq i_{k} \leq m}\left|\Lambda_{i_{k} i_{k}}\right|^{2} \geq \frac{3-\sqrt{5}}{2}\|\boldsymbol{\Lambda}\|^{2} .
\end{aligned}
$$

Therefore, there exists an index pair $\left(i_{k}, j_{k}\right)$ and $\nu_{k}(x)$ such that

$$
\frac{3}{3+\sqrt{5}} m(m-1)\left\|\partial \nu_{k}\left(\boldsymbol{I}_{2}\right)\right\|^{2} \geq\|\boldsymbol{\Lambda}\|^{2} .
$$

The proof is complete.
4.4. GU class of transformations. Inspired by the equation (13), we now propose the third class of elementary transformations: GU, which includes the plane upper triangular and plane diagonal transformations. In this case, we call Algorithm 1 the $B C D-G U$ algorithm, in which we choose a starting point $\boldsymbol{X}_{0} \in \mathbf{E U T}_{m}(\mathbb{C})$. Based on this class of elementary transformations, the subalgorithm to update $\boldsymbol{X}_{k}$ in Algorithm 1 is summarized in Subalgorithm 1c.

Let $\mathfrak{e u t}_{m}(\mathbb{C}) \subseteq \mathbb{C}^{m \times m}$ be the set of upper triangular matrices with trace equal to 0 . Then the tangent space to $\mathbf{E U T}_{m}(\mathbb{C})$ at a point $\boldsymbol{X} \in \mathbf{E U T}_{m}(\mathbb{C})$ can be constructed [7, 4] by $\mathbf{T}_{\boldsymbol{X}} \mathbf{E U T}_{m}(\mathbb{C})=\left\{\boldsymbol{X} \Omega, \Omega \in \mathfrak{e u t}_{m}(\mathbb{C})\right\}$ 。

Proposition 4.6. In Subalgorithm 1c, we can always choose an index pair $\left(i_{k}, j_{k}\right)$ and an elementary function $\nu_{k}(x)$ such that the inequality (38) is satisfied.

```
Subalgorithm 1c: The subalgorithm to update \(\boldsymbol{X}_{k}\) based on GU transformations
    Input: A fixed positive constant \(0<\varepsilon<\sqrt{\frac{2}{m(m+1)}}\).
    Output: New iterate \(\boldsymbol{X}_{k}\).
        Choose the index pair \(\left(i_{k}, j_{k}\right)\) and \(\nu_{k}\) satisfying (38), where \(\nu_{k}=\varphi_{k}\) or \(\rho_{k}\);
        Compute \(\boldsymbol{\Psi}_{k}^{*}\) that minimizes the function \(\nu_{k}\);
        Update \(\boldsymbol{X}_{k}=\boldsymbol{X}_{k-1} \boldsymbol{V}_{k}\), where \(\boldsymbol{V}_{k}=\boldsymbol{U}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}\) or \(\boldsymbol{D}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}\).
```

Proof. Note that $\boldsymbol{X}_{0} \in \mathbf{E U T}_{m}(\mathbb{C})$ in Subalgorithm 1c. We see that $\boldsymbol{X}_{k} \in \mathbf{E U T}_{m}(\mathbb{C})$ for all $k \in \mathbb{N}$. By Lemma 4.3, Lemma 4.5 and $\boldsymbol{\Lambda}=\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right) \in \mathfrak{e u t}_{m}(\mathbb{C})$, we get that

$$
\begin{aligned}
\sum_{1 \leq i_{k}<j_{k} \leq m}\left(\left\|\partial \varphi_{k}\left(\boldsymbol{I}_{2}\right)\right\|^{2}+\left\|\partial \rho_{k}\left(\boldsymbol{I}_{2}\right)\right\|^{2}\right) & =\sum_{1 \leq i_{k}<j_{k} \leq m}\left(\left|\Lambda_{i_{k} j_{k}}\right|^{2}+\left|\Lambda_{i_{k} i_{k}}-\Lambda_{j_{k} j_{k}}\right|^{2}\right) \\
& =\sum_{1 \leq i_{k}<j_{k} \leq m}\left|\Lambda_{i_{k} j_{k}}\right|^{2}+m \sum_{1 \leq i_{k} \leq m}\left|\Lambda_{i_{k} i_{k}}\right|^{2} \geq\|\boldsymbol{\Lambda}\|^{2} .
\end{aligned}
$$

Therefore, there exists an index pair $\left(i_{k}, j_{k}\right)$ and $\nu_{k}(x)$ such that

$$
\frac{m(m+1)}{2}\left\|\partial \nu_{k}\left(\boldsymbol{I}_{2}\right)\right\|^{2} \geq\|\boldsymbol{\Lambda}\|^{2}
$$

The proof is complete.

## 5. Triangular and diagonal transformations for Jadm problem

Let $f$ be the cost function (5). Let $\boldsymbol{\omega}_{k-1}=\left(\boldsymbol{U}_{k-1}, \boldsymbol{X}_{k-1}\right)$ and $g: \mathbf{S L}_{m}(\mathbb{C}) \rightarrow \mathbb{R}$ be the restricted function $f_{2, \boldsymbol{U}_{k-1}}$ as in Section 4.3. Denote $\boldsymbol{B}^{(\ell)}=\boldsymbol{U}_{k-1}^{\star} \boldsymbol{A}^{(\ell)} \boldsymbol{U}_{k-1}$. Then $g$ can be expressed as

$$
\begin{equation*}
g: \mathbf{S L}_{m}(\mathbb{C}) \rightarrow \mathbb{R}^{+}, \quad \boldsymbol{X} \mapsto \sum_{\ell=1}^{L} \mu_{\ell}\left\|\operatorname{offdiag}\left\{\boldsymbol{W}^{(\ell)}\right\}\right\|^{2}, \tag{39}
\end{equation*}
$$

where $\boldsymbol{W}^{(\ell)}=\boldsymbol{X}^{\star} \boldsymbol{B}^{(\ell)} \boldsymbol{X}$. In this section, we will first calculate the Riemannian gradient of $g$ in (39), and the partial derivatives of elementary functions $\varphi_{k}, \psi_{k}$ and $\rho_{k}$ in (37). Then, we will prove that conditions (26) and (27) are both satisfied in the triangular and diagonal transformations.
5.1. Riemannian gradient. Let $\boldsymbol{W}^{(\ell)}$ and $g$ be as in (39). Then, by (29) and (23), we have

$$
\begin{align*}
\nabla g(\boldsymbol{X}) & =2\left(\boldsymbol{X}^{\mathrm{H}}\right)^{-1} \sum_{\ell=1}^{L} \mu_{\ell} \Upsilon\left(\boldsymbol{W}^{(\ell)}\right),  \tag{40}\\
\operatorname{grad} g(\boldsymbol{X}) & =2 \boldsymbol{X} \sum_{\ell=1}^{L} \mu_{\ell}\left(\Upsilon\left(\boldsymbol{W}^{(\ell)}\right)-\frac{\operatorname{tr}\left(\Upsilon\left(\boldsymbol{W}^{(\ell)}\right)\right)}{n} \boldsymbol{I}_{n}\right) . \tag{41}
\end{align*}
$$

Remark 5.1. In the real case, the Euclidean gradient in (40) was earlier derived in [4, Eq. (6.3)] and [9, Section 2.3]. In this paper, we extend it to problem (39) in the complex case, and calculate the Riemannian gradieent (41) as well.
5.2. Elementary functions. Let $\boldsymbol{W}^{(\ell)}=\boldsymbol{X}_{k-1} \boldsymbol{B}^{(\ell)} \boldsymbol{X}_{k-1}$. Let

$$
\varrho=\left\{\begin{aligned}
1, & \text { if }(\cdot)^{\wedge}=(\cdot)^{\mathrm{H}} ; \\
-1, & \text { if }(\cdot)^{\star}=(\cdot)^{\mathrm{T}} .
\end{aligned}\right.
$$

Denote $(i, j)=\left(i_{k}, j_{k}\right)$ for simpliciy. Now we make the following notations:

$$
\begin{aligned}
& \text { - } \alpha_{1} \stackrel{\text { def }}{=} \sum_{\ell=1}^{L} \sum_{p \neq j} \mu_{\ell}\left(\left|\mathrm{W}_{i p}^{(\ell)}\right|^{2}+\left|\mathrm{W}_{p i}^{(\ell)}\right|^{2}\right), \\
& \alpha_{2} \stackrel{\text { def }}{=} \sum_{\ell=1}^{L} \sum_{p \neq j} \mu_{\ell}\left(\mathrm{W}_{i p}^{(\ell, \Re)} \mathrm{W}_{j p}^{(\ell, \Re)}+\mathrm{W}_{i p}^{(\ell, \Im)} \mathrm{W}_{j p}^{(\ell, \Im)}+\mathrm{W}_{p i}^{(\ell, \Re)} \mathrm{W}_{p j}^{(\ell, \Re)}+\mathrm{W}_{p i}^{(\ell, \Im)} \mathrm{W}_{p j}^{(\ell, \Im)}\right), \\
& \alpha_{3} \stackrel{\text { def }}{=} \sum_{\ell=1}^{L} \sum_{p \neq j} \mu_{\ell}\left(\varrho\left(\mathrm{W}_{i p}^{(\ell, \Im)} \mathrm{W}_{j p}^{(\ell, \Re)}-\mathrm{W}_{i p}^{(\ell, \Re)} \mathrm{W}_{j p}^{(\ell, \Im)}\right)+\mathrm{W}_{p i}^{(\ell, \Re)} \mathrm{W}_{p j}^{(\ell, \Im)}-\mathrm{W}_{p i}^{(\ell, \Im)} \mathrm{W}_{p j}^{(\ell, \Re)}\right) . \\
& \text { - } \beta_{1} \stackrel{\text { def }}{=} \sum_{\ell=1}^{L} \sum_{p \neq i} \mu_{\ell}\left(\left|\mathrm{W}_{j p}^{(\ell)}\right|^{2}+\left|\mathrm{W}_{p j}^{(\ell)}\right|^{2}\right), \\
& \beta_{2} \stackrel{\text { def }}{=} \sum_{\ell=1}^{L} \sum_{p \neq i} \mu_{\ell}\left(\mathrm{W}_{i p}^{(\ell, \Re)} \mathrm{W}_{j p}^{(\ell, \Re)}+\mathrm{W}_{i p}^{(\ell, \Im)} \mathrm{W}_{j p}^{(\ell, \Im)}+\mathrm{W}_{p i}^{(\ell, \Re)} \mathrm{W}_{p j}^{(\ell, \Re)}+\mathrm{W}_{p i}^{(\ell, \Im)} \mathrm{W}_{p j}^{(\ell, \Im)}\right), \\
& \beta_{3} \stackrel{\text { def }}{=} \sum_{\ell=1}^{L} \sum_{p \neq i} \mu_{\ell}\left(\varrho\left(\mathrm{W}_{i p}^{(\ell, \Re)} \mathrm{W}_{j p}^{(\ell, \Im)}-\mathrm{W}_{i p}^{(\ell, \Im)} \mathrm{W}_{j p}^{(\ell, \Re)}\right)+\mathrm{W}_{p i}^{(\ell, \Im)} \mathrm{W}_{p j}^{(\ell \Re)}-\mathrm{W}_{p i}^{(\ell, \Re)} \mathrm{W}_{p j}^{(\ell, \Im)}\right) . \\
& \\
& \\
& \gamma_{1} \stackrel{\text { def }}{=} \sum_{\ell=1}^{L} \sum_{p \neq i, j} \mu_{\ell}\left(\left|\mathrm{W}_{i p}^{(\ell)}\right|^{2}+\left|\mathrm{W}_{p i}^{(\ell)}\right|^{2}\right), \quad \gamma_{2} \stackrel{\text { def }}{=} \sum_{\ell=1}^{L} \sum_{p \neq i, j} \mu_{\ell}\left(\left|\mathrm{W}_{j p}^{(\ell)}\right|^{2}+\left|\mathrm{W}_{p j}^{(\ell)}\right|^{2}\right) .
\end{aligned}
$$

Then we can get the following results by direct calculations.
Lemma 5.2. Let the function $g$ be as in (39). Then
(i) the elementary function $\varphi_{k}$ in (37) satisfies

$$
\begin{align*}
\varphi_{k}(x, y)-\varphi_{k}(0,0) & =\alpha_{1} x^{2}+2 \alpha_{2} x+\alpha_{1} y^{2}+2 \alpha_{3} y \\
\varphi_{k}\left(x_{k}^{*}, y_{k}^{*}\right)-\varphi_{k}(0,0) & =-\frac{1}{\alpha_{1}}\left(\alpha_{2}^{2}+\alpha_{3}^{2}\right)  \tag{42}\\
\partial \varphi_{k}(0,0) & =2\left[\alpha_{2}, \alpha_{3}\right]^{\top}
\end{align*}
$$

(ii) the elementary function $\psi_{k}$ in (37) satisfies

$$
\begin{aligned}
\psi_{k}(x, y)-\psi_{k}(0,0) & =\beta_{1} x^{2}+2 \beta_{2} x+\beta_{1} y^{2}+2 \beta_{3} y \\
\psi_{k}\left(x_{k}^{*}, y_{k}^{*}\right)-\psi_{k}(0,0) & =-\frac{1}{\beta_{1}}\left(\beta_{2}^{2}+\beta_{3}^{2}\right) \\
\partial \psi_{k}(0,0) & =2\left[\beta_{2}, \beta_{3}\right]^{\top}
\end{aligned}
$$

(iii) the elementary function $\rho_{k}$ in (37) satisfies

$$
\begin{aligned}
\rho_{k}(x, y)-\rho_{k}(1,0) & =\gamma_{1}\left(x^{2}+y^{2}\right)+\gamma_{2} \frac{1}{x^{2}+y^{2}}-\gamma_{1}-\gamma_{2}, \\
\rho_{k}\left(x_{k}^{*}, y_{k}^{*}\right)-\rho_{k}(1,0) & =-\left(\sqrt{\gamma_{1}}-\sqrt{\gamma_{2}}\right)^{2}, \\
\partial \rho_{k}(1,0) & =2\left[\gamma_{1}-\gamma_{2}, 0\right]^{\top} .
\end{aligned}
$$

Remark 5.3. In the complex case, the solution $z_{k}^{*}=x_{k}^{*}+i y_{k}^{*}$ in (42) was earlier derived in [45, Eq. (8)]. In the real case, the solution $x_{k}^{*}$ in (42) was earlier derived in [5, Eq. (7)].

Setting 5.4. In Algorithm 1 for cost function (5), when $\nu_{k}=\varphi_{k}$, we see that $x_{k}^{*}=0$ if $\alpha_{1} \neq 0$ and $\alpha_{2}=0$. It is not possible that $\alpha_{1}=0$ and $\alpha_{2} \neq 0$. If $\alpha_{1}=\alpha_{2}=0$, we set $x_{k}^{*}=0$. In the case of $\nu_{k}=\psi_{k}$, we make the similar settings for the value of $y_{k}^{*}$.

Setting 5.5. Let $0<\varsigma<\frac{1}{4}$ be a small positive constant. In Algorithm 1 for cost function (5), if $\nu_{k}=\rho_{k}$, we always set $y_{k}^{*}=0$. Moreover, we determine $x_{k}^{*}$ based on the following rules.

- If $\gamma_{1}=\gamma_{2}=0$, we set $x_{k}^{*}=0$.
- Let $\varpi \stackrel{\text { def }}{=} \frac{\gamma_{2}}{\gamma_{1}}$. If $\varpi \in[0, \varsigma)$, we set $x_{k}^{*}=\frac{1}{2}$. If $\varpi \in\left(\frac{1}{\varsigma},+\infty\right]$, we set $x_{k}^{*}=2$.
- Otherwise, if $\varpi \in\left[\varsigma, \frac{1}{\varsigma}\right]$, we set $x_{k}^{*}=\sqrt[4]{\varpi}$, which is the minimum point.
5.3. Condition (26) for global convergence. It will be seen that $f\left(\boldsymbol{\omega}_{k}\right) \leq f\left(\boldsymbol{\omega}_{k-1}\right)$ always holds in Algorithm 1. In Algorithm 1 for cost function (5), we denote $\mathrm{M}_{0} \stackrel{\text { def }}{=} f\left(\boldsymbol{\omega}_{0}\right)$. Then we have that $\gamma_{1}+\gamma_{2} \leq \mathrm{M}_{0}=f\left(\boldsymbol{\omega}_{0}\right)$.

Lemma 5.6. In Algorithm 1 for cost function (5), there exists $\iota_{\rho}>0$ such that

$$
\begin{equation*}
g\left(\boldsymbol{X}_{k-1}\right)-g\left(\boldsymbol{X}_{k}\right) \geq \iota_{\rho}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\|, \tag{43}
\end{equation*}
$$

whenever $\nu_{k}=\rho_{k}$.
Proof. We now prove the inequality (43) by Lemma 5.2 (iii) in three different cases shown in Setting 5.5.

- If $\gamma_{1}=\gamma_{2}=0$, it is clear that the inequality (43) is satisfied for any $\iota_{\rho}>0$.
- If $\varpi \in[0, \varsigma)$, we get that

$$
\begin{aligned}
\rho_{k}(1,0)-\rho_{k}\left(x_{k}^{*}, y_{k}^{*}\right) & =\frac{3}{2} \gamma_{1}(1-4 \varpi)=\frac{3(1-4 \varpi)}{8(1-\varpi)}\left|\partial \rho_{k}(1,0)\right| \geq \frac{3(1-4 \varpi) \varepsilon}{8(1-\varpi)}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\| \\
& \geq \frac{3(1-4 \varpi) \varepsilon}{4 \sqrt{5}(1-\varpi)} \frac{\sqrt{5}}{2}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\| \geq \frac{3(1-4 \varpi) \varepsilon}{4 \sqrt{5}(1-\varpi)}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\| \\
& \geq \frac{3(1-4 \varsigma) \varepsilon}{4 \sqrt{5}}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\| .
\end{aligned}
$$

- If $\varpi \in\left(\frac{1}{\varsigma},+\infty\right]$, we similarly get that

$$
\begin{aligned}
\rho_{k}(1,0)-\rho_{k}\left(x_{k}^{*}, y_{k}^{*}\right) & =\frac{3}{2} \gamma_{2}\left(1-\frac{4}{\varpi}\right)=\frac{3\left(1-\frac{4}{\omega}\right)}{8\left(1-\frac{1}{\omega}\right)}\left|\partial \rho_{k}(1,0)\right| \geq \frac{3\left(1-\frac{4}{\omega}\right) \varepsilon}{8\left(1-\frac{1}{\omega}\right)}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\| \\
& \geq \frac{3\left(1-\frac{4}{\omega}\right) \varepsilon}{4 \sqrt{5}\left(1-\frac{1}{\omega}\right)} \frac{\sqrt{5}}{2}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\| \geq \frac{3\left(1-\frac{4}{\omega}\right) \varepsilon}{4 \sqrt{5}\left(1-\frac{1}{\omega}\right)}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\| \\
& \geq \frac{3(1-4 \varsigma) \varepsilon}{4 \sqrt{5}}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\| .
\end{aligned}
$$

- If $\varpi \in\left[\varsigma, \frac{1}{\varsigma}\right]$, it is easy to verify that

$$
\begin{equation*}
\sqrt[4]{\gamma_{1} \gamma_{2}} \geq \frac{\sqrt[4]{\varsigma}}{2}\left(\sqrt{\gamma_{1}}+\sqrt{\gamma_{2}}\right) \tag{44}
\end{equation*}
$$

Then, we get that

$$
\begin{aligned}
\rho_{k}(1,0)-\rho_{k}\left(x_{k}^{*}, y_{k}^{*}\right) & =2\left(\sqrt{\gamma_{1}}-\sqrt{\gamma_{2}}\right)^{2}=\frac{1}{2}\left|\partial \rho_{k}(1,0)\right| \frac{\left|\sqrt{\gamma_{1}}-\sqrt{\gamma_{2}}\right|}{\sqrt{\gamma_{1}}+\sqrt{\gamma_{2}}} \\
& \geq \frac{\varepsilon \sqrt[4]{\varsigma}}{4}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\| \frac{\left|\sqrt{\gamma_{1}}-\sqrt{\gamma_{2}}\right|}{\sqrt[4]{\gamma_{1} \gamma_{2}}} \quad(\text { by equation (44)) } \\
& \geq \frac{\varepsilon \sqrt[4]{\varsigma}}{4}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\| \frac{\left|\sqrt[4]{\gamma_{1}}-\sqrt[4]{\gamma_{2}}\right|\left(\sqrt{\gamma_{1}}+\sqrt{\gamma_{2}}\right)^{1 / 2}}{\sqrt[4]{\gamma_{1} \gamma_{2}}} \\
& \geq \frac{\varepsilon \sqrt[4]{\varsigma}}{4}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|\left|x_{k}^{*}\right| \sqrt{1+\frac{1}{\left(1+x_{k}^{*}\right)^{2}}} \geq \frac{\varepsilon \sqrt[4]{\varsigma}}{4}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\| .
\end{aligned}
$$

Now we set $\iota_{\rho}=\min \left(\frac{3(1-4 \varsigma) \varepsilon}{4 \sqrt{5}}, \frac{\varepsilon \sqrt[4]{5}}{4}\right)$. The proof is complete.
Note that $\left\|\boldsymbol{X}_{k}-\boldsymbol{X}_{k-1}\right\| \leq\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\|\left\|\boldsymbol{X}_{k-1}\right\|$ and $\left\|\operatorname{grad} g\left(\boldsymbol{X}_{k-1}\right)\right\| \leq\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|\left\|\boldsymbol{X}_{k-1}\right\|$. Let $\iota_{\rho}$ be as in (43) and $\sigma_{\rho}=\frac{\iota_{\rho}}{\mathrm{M}_{\omega}^{2}}>0$. Then the next result follows directly from Lemma 5.6.

Corollary 5.7. In Algorithm 1 for cost function (5), if the condition (17) is satisfied, then

$$
\begin{equation*}
g\left(\boldsymbol{X}_{k-1}\right)-g\left(\boldsymbol{X}_{k}\right) \geq \sigma_{\rho}\left\|\operatorname{grad} g\left(\boldsymbol{X}_{k-1}\right)\right\|\left\|\boldsymbol{X}_{k}-\boldsymbol{X}_{k-1}\right\|, \tag{45}
\end{equation*}
$$

whenever $\nu_{k}=\rho_{k}$.
Lemma 5.8. In Algorithm 1 for cost function (5), there exists $\iota_{\varphi}>0$ such that

$$
\begin{equation*}
g\left(\boldsymbol{X}_{k-1}\right)-g\left(\boldsymbol{X}_{k}\right) \geq \iota_{\varphi}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\| \tag{46}
\end{equation*}
$$

whenever $\nu_{k}=\varphi_{k}$ or $\psi_{k}$.
Proof. We prove that the inequality (46) is satisfied in two cases.

- If $\nu_{k}=\varphi_{k}$, by Lemma 5.2(i), we see that

$$
\begin{aligned}
g\left(\boldsymbol{X}_{k-1}\right)-g\left(\boldsymbol{X}_{k}\right) & =\varphi_{k}(0,0)-\varphi_{k}\left(x_{k}^{*}, y_{k}^{*}\right)=2\left(\frac{\left(\alpha_{2}\right)^{2}}{\alpha_{1}}+\frac{\left(\alpha_{3}\right)^{2}}{\alpha_{1}}\right) \\
& =\frac{1}{2}\left\|\partial \varphi_{k}(0,0)\right\|\left\|\left(x_{k}^{*}, y_{k}^{*}\right)\right\| \geq \frac{\varepsilon}{2}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\| .
\end{aligned}
$$

- If $\nu_{k}=\psi_{k}$, by Lemma 5.2(ii), we see that

$$
\begin{aligned}
g\left(\boldsymbol{X}_{k-1}\right)-g\left(\boldsymbol{X}_{k}\right) & =\psi_{k}(0,0)-\psi_{k}\left(x_{k}^{*}, y_{k}^{*}\right)=2\left(\frac{\left(\beta_{2}\right)^{2}}{\beta_{1}}+\frac{\left(\beta_{3}\right)^{2}}{\beta_{1}}\right) \\
& =\frac{1}{2}\left\|\partial \psi_{k}(0,0)\right\|\left\|\left(x_{k}^{*}, y_{k}^{*}\right)\right\| \geq \frac{\varepsilon}{2}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\|
\end{aligned}
$$

Now we set $\iota_{\varphi}=\frac{\varepsilon}{2}$. The proof is complete.
Let $\iota_{\varphi}$ be as in (46) and $\sigma_{\varphi}=\frac{\iota_{\varphi}}{\mathrm{M}_{\omega}^{2}}>0$. Similar as for Corollary 5.7, the next result follows directly from Lemma 5.8.

Corollary 5.9. In Algorithm 1 for cost function (5), if the condition (17) is satisfied, then

$$
g\left(\boldsymbol{X}_{k-1}\right)-g\left(\boldsymbol{X}_{k}\right) \geq \sigma_{\varphi}\left\|\operatorname{grad} g\left(\boldsymbol{X}_{k-1}\right)\right\|\left\|\boldsymbol{X}_{k}-\boldsymbol{X}_{k-1}\right\|
$$

whenever $\nu_{k}=\varphi_{k}$ or $\psi_{k}$.

### 5.4. Condition (27) for weak convergence.

Lemma 5.10. In Algorithm 1 for cost function (5), if the condition (17) is satisfied, then there exists $\kappa_{\rho}>0$ such that

$$
\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\| \geq \kappa_{\rho}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|
$$

whenever $\nu_{k}=\rho_{k}, \varphi_{k}$ or $\psi_{k}$.
Proof. If $\nu_{k}=\varphi_{k}$, we have

$$
\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\|^{2}=\frac{\alpha_{2}^{2}+\alpha_{3}^{2}}{\alpha_{1}^{2}} \geq \frac{\left\|\partial \varphi_{k}(0,0)\right\|^{2}}{4 \mathrm{M}_{2}^{2}} \geq \frac{\varepsilon^{2}}{4 \mathrm{M}_{2}^{2}}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|^{2}
$$

The case $\nu_{k}=\psi_{k}$ is similar. Now we prove the $\nu_{k}=\rho_{k}$ case. If $\varpi \in[0, \varsigma)$, we have

$$
\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\|^{2}=\frac{5}{4} \geq \frac{5}{4} \frac{1}{8\left(\gamma_{1}^{2}+\gamma_{2}^{2}\right)}\left\|\partial \rho_{k}(1,0)\right\|^{2} \geq \frac{5}{4} \frac{\varepsilon^{2}}{16 \mathrm{M}_{\omega}^{2}}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|^{2}
$$

The case $\varpi \in\left(\frac{1}{\varsigma},+\infty\right]$ is similar. If $\varpi \in\left[\varsigma, \frac{1}{\varsigma}\right]$, we have

$$
\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\|^{2} \geq\left(1-x_{k}^{*}\right)^{2} \geq \frac{1}{\mathrm{M}_{\omega}^{2} \mathrm{M}_{3}^{2}} \gamma_{1}^{2}(1-\varpi)^{2}=\frac{1}{4 \mathrm{M}_{\omega}^{2} \mathrm{M}_{3}^{2}}\left\|\partial \rho_{k}(1,0)\right\|^{2} \geq \frac{\varepsilon^{2}}{4 \mathrm{M}_{\omega}^{2} \mathrm{M}_{3}^{2}}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|^{2}
$$

We only need to set $\kappa_{\rho}^{2}=\min \left(\frac{\varepsilon^{2}}{4 \mathrm{M}_{2}^{2}}, \frac{5}{4} \frac{\varepsilon^{2}}{16 \mathrm{M}_{\omega}^{2}}, \frac{\varepsilon^{2}}{4 \mathrm{M}_{\omega}^{2} \mathrm{M}_{3}^{2}}\right)$. The proof is complete.
By Lemma 5.6, Lemma 5.8 and Lemma 5.10, we can easily get the following results by setting $\eta_{\rho}=\sigma_{\rho} \kappa_{\rho}$ and $\eta_{\varphi}=\sigma_{\varphi} \kappa_{\rho}$.
Corollary 5.11. In Algorithm 1 for cost function (5), if the condition (17) is satisfied, then

$$
\begin{equation*}
g\left(\boldsymbol{X}_{k-1}\right)-g\left(\boldsymbol{X}_{k}\right) \geq \eta_{\rho}\left\|\operatorname{grad} g\left(\boldsymbol{X}_{k-1}\right)\right\|^{2} \tag{47}
\end{equation*}
$$

whenever $\nu_{k}=\rho_{k}$.
Corollary 5.12. In Algorithm 1 for cost function (5), if the condition (17) is satisfied, then

$$
g\left(\boldsymbol{X}_{k-1}\right)-g\left(\boldsymbol{X}_{k}\right) \geq \eta_{\varphi}\left\|\operatorname{grad} g\left(\boldsymbol{X}_{k-1}\right)\right\|^{2}
$$

whenever $\nu_{k}=\varphi_{k}$ or $\psi_{k}$.

## 6. Givens plane transformations for JADM problem

Let $g$ be as in (39). Let $\boldsymbol{W}^{(\ell)}$ and $\varrho$ be as in Section 5.2. Denote $(i, j)=\left(i_{k}, j_{k}\right)$ for simpliciy. Define

$$
\boldsymbol{\Gamma}^{\left(i, j, \boldsymbol{X}_{k-1}\right)} \stackrel{\text { def }}{=} \frac{\varrho}{2} \sum_{\ell=1}^{L} \mu_{\ell} \Re\left(\boldsymbol{z}_{i, j}\left(\boldsymbol{W}^{(\ell)}\right) \boldsymbol{z}_{i, j}^{\mathrm{H}}\left(\boldsymbol{W}^{(\ell)}\right)\right)
$$

where

$$
\boldsymbol{z}_{i, j}(\boldsymbol{W}) \stackrel{\text { def }}{=}\left\{\begin{array}{lll}
{\left[\begin{array}{lll}
\mathrm{W}_{j j}-\mathrm{W}_{i i} & \mathrm{~W}_{i j}+\mathrm{W}_{j i} & -i\left(\mathrm{~W}_{i j}-\mathrm{W}_{j i}\right)
\end{array}\right]^{\top},} & \text { if }(\cdot)^{\wedge}=(\cdot)^{\mathrm{H}} \\
{\left[\begin{array}{lll}
\mathrm{W}_{i j}+\mathrm{W}_{j i} & \mathrm{~W}_{i i}-\mathrm{W}_{j j} & i\left(\mathrm{~W}_{i i}+\mathrm{W}_{j j}\right)
\end{array}\right]^{\top},} & \text { if }(\cdot)^{\star}=(\cdot)^{\top}
\end{array}\right.
$$

Let

$$
c_{0}= \begin{cases}\frac{1}{2} \sum_{\ell=1}^{L} \mu_{\ell}\left|\mathrm{W}_{j j}^{(\ell)}-\mathrm{W}_{i i}^{(\ell)}\right|^{2}, & \text { if }(\cdot)^{\star}=(\cdot)^{\mathrm{H}} ; \\ -\frac{1}{2} \sum_{\ell=1}^{L} \mu_{\ell}\left|\mathrm{W}_{i j}^{(\ell)}+\mathrm{W}_{j i}^{(\ell)}\right|^{2}, & \text { if }(\cdot)^{\star}=(\cdot)^{\top}\end{cases}
$$

6.1. Elementary function. As in [44, Eq. (4.4)], we let

$$
\begin{equation*}
\boldsymbol{r}=\boldsymbol{r}\left(c, s_{1}, s_{2}\right) \stackrel{\text { def }}{=}\left[2 c^{2}-1,-2 c s_{1},-2 c s_{2}\right]^{\top}=[\cos 2 \theta,-\sin 2 \theta \cos \phi,-\sin 2 \theta \sin \phi]^{\top} \tag{48}
\end{equation*}
$$

Then we can get the following result ${ }^{11}$ by direct calculations.
Lemma 6.1. In Algorithm 1 for cost function (5), the elementary function $h_{k}$ satisfies

$$
\begin{equation*}
h_{k}\left(c, s_{1}, s_{2}\right)-h_{k}(1,0,0)=-\left(\boldsymbol{r}^{\top} \boldsymbol{\Gamma}^{\left(i, j, \boldsymbol{X}_{k-1}\right)} \boldsymbol{r}-c_{0}\right) \tag{49}
\end{equation*}
$$

Now we denote $\boldsymbol{\Gamma}=\boldsymbol{\Gamma}^{\left(i, j, \boldsymbol{X}_{k-1}\right)}$ for simplicity. It follows by (48) and (49) that

$$
\begin{equation*}
h_{k}\left(c, s_{1}, s_{2}\right)-h_{k}(1,0,0)=-\left(q(\theta, \phi)-c_{0}\right) \tag{50}
\end{equation*}
$$

where

$$
\begin{align*}
q(\theta, \phi) \stackrel{\text { def }}{=} & \frac{1}{2}\left(\Gamma_{11}-\Gamma_{22} \cos ^{2} \phi-\Gamma_{33} \sin ^{2} \phi-\Gamma_{23} \sin (2 \phi)\right) \cos (4 \theta) \\
& -\left(\Gamma_{12} \cos \phi+\Gamma_{13} \sin \phi\right) \sin (4 \theta)+\frac{1}{2}\left(\Gamma_{11}+\Gamma_{22} \cos ^{2} \phi+\Gamma_{33} \sin ^{2} \phi+\Gamma_{23} \sin (2 \phi)\right) \tag{51}
\end{align*}
$$

Note that, by Lemma 6.1 and (48), we have

$$
\partial h_{k}\left(\boldsymbol{I}_{2}\right)=-4\left[\begin{array}{lll}
0 & \Gamma_{12} & \Gamma_{13} \tag{52}
\end{array}\right]^{\top}
$$

Remark 6.2. By equation (50), we see that $h_{k}\left(\theta+\frac{\pi}{2}, \phi\right)=h_{k}(\theta, \phi)$ for any $\theta, \phi \in \mathbb{R}$. Therefore, we can always choose $\theta_{*} \in\left[-\frac{\pi}{4}, \frac{\pi}{4}\right]$.
Setting 6.3. In Algorithm 1 for cost function (5), we set a positive constant $\epsilon>0$. If $\nu_{k}=h_{k}$, we find the eigenvector $\boldsymbol{u}$ of $\boldsymbol{\Gamma}$ corresponding to the largest eigenvalue. Define two vectors $\boldsymbol{v}_{i, j} \stackrel{\text { def }}{=}\left[\begin{array}{ll}\Gamma_{12} & \Gamma_{13}\end{array}\right]^{\top} \in \mathbb{R}^{2}$ and $\boldsymbol{w}_{i, j} \stackrel{\text { def }}{=}\left[\begin{array}{ll}u_{2} & u_{3}\end{array}\right]^{\top} \in \mathbb{R}^{2}$.

- If it holds that

$$
\begin{equation*}
\left|\left\langle\boldsymbol{v}_{i, j}, \boldsymbol{w}_{i, j}\right\rangle\right| \geq \epsilon\left\|\boldsymbol{v}_{i, j}\right\|\left\|\boldsymbol{w}_{i, j}\right\| \tag{53}
\end{equation*}
$$

then we find $\phi_{*}$ and $\theta_{*}$ by setting $\boldsymbol{r}=\boldsymbol{u}$, and $\boldsymbol{\Psi}_{k}^{*}=\boldsymbol{\Psi}\left(\theta_{*}, \phi_{*}\right)$;

- Otherwise, we set $\left[\cos \phi_{*} \sin \phi_{*}\right]^{\top}=\frac{\boldsymbol{v}_{i, j}}{\left\|\boldsymbol{v}_{i, j}\right\|}$ and then determine the best $\theta_{*}$ based on $\phi_{*}$.

[^8]6.2. Condition (26) for global convergence. We first present a lemma, which will help us to prove Lemma 6.5.

Lemma 6.4. Let $\alpha, \beta \in \mathbb{R}$ be two constants. For $\theta \in\left[-\frac{\pi}{4}, \frac{\pi}{4}\right]$, we define

$$
p(\theta) \stackrel{\text { def }}{=} \alpha \cos (4 \theta)+\beta \sin (4 \theta) .
$$

If $p\left(\theta_{*}\right)=\max p(\theta)$, we have

$$
p\left(\theta_{*}\right)-p(0) \geq 2 \sqrt{2}|\beta|\left|\sin \left(\frac{\theta_{*}}{2}\right)\right| .
$$

Lemma 6.5. Let the function $q(\theta, \phi)$ be as in (51). Suppose that $\phi_{*}$ and $\theta_{*}$ are determined as in Setting 6.3. Then we have

$$
q\left(\theta_{*}, \phi_{*}\right)-q(0,0) \geq 2 \sqrt{2} \epsilon\left|\sin \left(\frac{\theta_{*}}{2}\right)\right|\left\|\boldsymbol{v}_{i, j}\right\| .
$$

Proof. By Setting 6.3, we see that

$$
\begin{equation*}
\left|\left\langle\boldsymbol{v}_{i, j},\left[\cos \phi_{*} \sin \phi_{*}\right]^{\top}\right\rangle\right| \geq \epsilon\left\|\boldsymbol{v}_{i, j}\right\| \tag{54}
\end{equation*}
$$

always holds. By Lemma 6.4 and (54), we get that

$$
\begin{aligned}
q\left(\theta_{*}, \phi_{*}\right)-q(0,0) & =q\left(\theta_{*}, \phi_{*}\right)-q\left(0, \phi_{*}\right) \geq 2 \sqrt{2}\left|\sin \left(\frac{\theta_{*}}{2}\right)\right|\left|\Gamma_{12} \cos \phi_{*}+\Gamma_{13} \sin \phi_{*}\right| \\
& \geq 2 \sqrt{2} \epsilon\left|\sin \left(\frac{\theta_{*}}{2}\right)\right|\left\|\boldsymbol{v}_{i, j}\right\| .
\end{aligned}
$$

The proof is complete.
Lemma 6.6. In Algorithm 1 for cost function (5), there exists $\iota_{h}>0$ such that

$$
\begin{equation*}
g\left(\boldsymbol{X}_{k-1}\right)-g\left(\boldsymbol{X}_{k}\right) \geq \iota_{h}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\|, \tag{55}
\end{equation*}
$$

whenever $\nu_{k}=h_{k}$.
Proof. We only prove the $(\cdot)^{\boldsymbol{\wedge}}=(\cdot)^{\mathrm{H}}$ case. The other case is similar. By Lemma 6.5 and (52), we get that

$$
\begin{aligned}
h_{k}(0,0)-h_{k}\left(\theta_{*}, \phi_{*}\right) & \left.\geq 2 \sqrt{2} \epsilon\left|\sin \left(\frac{\theta_{*}}{2}\right)\right|\left\|\boldsymbol{v}_{i, j}\right\|=\frac{\epsilon}{4} 2 \sqrt{2} \right\rvert\, \sin \left(\frac{\theta_{*}}{2}\right)\left\|\partial h_{k}\left(\boldsymbol{I}_{2}\right)\right\| \\
& \geq \frac{\epsilon \varepsilon}{4}\left\|\boldsymbol{G}^{\left(i, j, \boldsymbol{\Psi}_{k}^{*}\right)}-\boldsymbol{I}_{m}\right\|\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\| .
\end{aligned}
$$

We can set $\iota_{h}=\frac{\epsilon \varepsilon}{4}$. The proof is complete.
Let $\iota_{h}$ be as in (55) and $\sigma_{h}=\frac{\iota_{h}}{\mathrm{M}_{\omega}^{2}}>0$. Similar as for Corollary 5.7, the next result follows directly from Lemma 6.6.

Corollary 6.7. In Algorithm 1 for cost function (5), if the condition (17) is always satisfied, then

$$
g\left(\boldsymbol{X}_{k-1}\right)-g\left(\boldsymbol{X}_{k}\right) \geq \sigma_{h}\left\|\operatorname{grad} g\left(\boldsymbol{X}_{k-1}\right)\right\|\left\|\boldsymbol{X}_{k}-\boldsymbol{X}_{k-1}\right\|,
$$

whenever $\nu_{k}=h_{k}$.
6.3. Condition (27) for weak convergence. If the condition (17) is satisfied, it is easy to see that there exists $\mathrm{M}_{\Gamma}>0$ such that $\left\|\boldsymbol{\Gamma}^{\left(i, j, \boldsymbol{X}_{k-1}\right)}\right\| \leq \mathrm{M}_{\Gamma}$ always hods.

Lemma 6.8. In Algorithm 1 for cost function (5), if the condition (17) is satisfied, then there exists $\kappa_{h}>0$ such that

$$
\begin{equation*}
\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\| \geq \kappa_{h}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\|, \tag{56}
\end{equation*}
$$

whenever $\nu_{k}=h_{k}$.
Proof. If $\boldsymbol{v}_{i, j}$ and $\boldsymbol{w}_{i, j}$ satisfy the condition (53), then the inequality (56) can be proved by a similar method as for [44, Lemma 7.2]. Otherwise, if we set $\left[\cos \phi_{*} \sin \phi_{*}\right]^{\top}=\frac{\boldsymbol{v}_{i, j}}{\left\|\boldsymbol{v}_{i, j}\right\|}$ and find $\theta_{*}$ based on $\phi_{*}$, then

$$
\begin{aligned}
\left|\sin \left(4 \theta_{*}\right)\right| & =\frac{\left|\Gamma_{12} \cos \phi_{*}+\Gamma_{13} \sin \phi_{*}\right|}{\sqrt{\left(\Gamma_{12} \cos \phi_{*}+\Gamma_{13} \sin \phi_{*}\right)^{2}+\frac{1}{4}\left(\Gamma_{11}-\Gamma_{22} \cos ^{2} \phi_{*}-\Gamma_{33} \sin ^{2} \phi_{*}-\Gamma_{23} \sin \left(2 \phi_{*}\right)\right)^{2}}} \\
& \geq \frac{\sqrt{\Gamma_{12}^{2}+\Gamma_{13}^{2}}}{2 \sqrt{5} \mathrm{M}_{\Gamma}}=\frac{\left\|\partial h_{k}\left(\boldsymbol{I}_{2}\right)\right\|}{8 \sqrt{5} \mathrm{M}_{\Gamma}} \geq \frac{\varepsilon}{8 \sqrt{5} \mathrm{M}_{\Gamma}}\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\| .
\end{aligned}
$$

Note that

$$
\left\|\boldsymbol{\Psi}_{k}^{*}-\boldsymbol{I}_{2}\right\|=2 \sqrt{2}\left|\sin \left(\frac{\theta_{*}}{2}\right)\right| \geq \frac{\sqrt{2}}{4}\left|\sin \left(4 \theta_{*}\right)\right| .
$$

We only need to set $\kappa_{h}=\frac{\sqrt{2} \varepsilon}{32 \sqrt{5} \mathrm{M}_{\Gamma}}$ in this case. The proof is complete.
Then, by Lemma 6.6 and Lemma 6.8, we can easily get the following result by setting $\eta_{h}=\sigma_{h} \kappa_{h}$.

Corollary 6.9. In Algorithm 1 for cost function (5), if the condition (17) is satisfied, then

$$
g\left(\boldsymbol{X}_{k-1}\right)-g\left(\boldsymbol{X}_{k}\right) \geq \eta_{h}\left\|\operatorname{grad} g\left(\boldsymbol{X}_{k-1}\right)\right\|^{2}
$$

whenever $\nu_{k}=h_{k}$.

## 7. Convergence analysis

### 7.1. Convergence analysis of BCD-GLU and BCD-GQU algorithms.

7.1.1. Global convergence. We first prove the following result about the global convergence of BCD-GLU and BCD-GQU algorithms for cost function (5).

Theorem 7.1. In BCD-GLU and BCD-GQU algorithms for cost function (5), if the condition (17) is satisfied, then the iterates $\left\{\boldsymbol{\omega}_{k}\right\}_{k \geq 1}$ produced by BCD-GLU and BCD-GQU algorithms converge to a point $\boldsymbol{\omega}_{*}$.

Proof. By Lemma 3.3, Corollary 5.7, Corollary 5.9 and Corollary 6.7, we see that the condition (26) is always satisfied in BCD-GLU and BCD-GQU algorithms for cost function (5), if we set $\sigma=\min \left(\sigma_{p}, \sigma_{\rho}, \sigma_{\varphi}, \sigma_{h}\right)$. Then the proof is complete by Proposition 2.4.
7.1.2. Weak convergence. Now we prove the following result about the weak convergence of BCD-GLU and BCD-GQU algorithms for cost function (5).
Theorem 7.2. In BCD-GLU and BCD-GQU algorithms for cost function (5), if condition (17) is always satisfied, and $\boldsymbol{\omega}_{*}$ is an accumulation point of the iterates $\left\{\boldsymbol{\omega}_{k}\right\}_{k \geq 1}$, then $\boldsymbol{\omega}_{*}$ is a stationary point of the cost function (5).

Proof. By Equation (35), Corollary 5.11, Corollary 5.12 and Corollary 6.9, we see that the condition (27) is always satisfied in BCD-GLU and BCD-GQU algorithms for cost function (5), if we set $\eta=\min \left(\eta_{p}, \eta_{\rho}, \eta_{\varphi}, \eta_{h}\right)$. Then the proof is complete by Lemma 2.5.

### 7.2. Convergence analysis of BCD-GU algorithm.

7.2.1. Global convergence. Similar as in Section 7.1, we have the following result about the global convergence of BCD-GU algorithm for cost function (5).

Theorem 7.3. In BCD-GU algorithm for cost function (5), if the condition (17) is satisfied, then the iterates $\left\{\boldsymbol{\omega}_{k}\right\}_{k \geq 1}$ produced by BCD-GU algorithm converge to a point $\boldsymbol{\omega}_{*}$.
7.2.2. Weak convergence. We also have the following result about the weak convergence of BCD-GU algorithm for cost function (5).
Theorem 7.4. In BCD-GU algorithm for cost function (5), if condition (17) is always satisfied, and $\boldsymbol{\omega}_{*}$ is an accumulation point of the iterates $\left\{\boldsymbol{\omega}_{k}\right\}_{k \geq 1}$, then $\boldsymbol{\omega}_{*}$ is a stationary point of the cost function (5).

### 7.3. Jacobi-type algorithms and their convergence analysis.

7.3.1. Jacobi-GLU and Jacobi-GQU algorithms on $\mathbf{S L}_{m}(\mathbb{C})$. Let $n=m$ in the cost function (8) and fix $\boldsymbol{U}_{k}=\boldsymbol{U}_{0}$ for $k \geq 1$ in Algorithm 1. In other words, we keep the first block variable unchanged and only update the second block variable $\boldsymbol{X}_{k}$ in $\mathbf{S L}_{m}(\mathbb{C})$ using GLU and GQU transformations ${ }^{12}$ in Section 4.3. Then we get the Jacobi-GLU and Jacobi-GQU algorithms ${ }^{13}$ to minimize the restricted function

$$
g: \mathbf{S L}_{m}(\mathbb{C}) \rightarrow \mathbb{R}^{+}, \quad \boldsymbol{X} \mapsto f\left(\boldsymbol{U}_{0}, \boldsymbol{X}\right)
$$

7.3.2. Special cases: Jacobi-GLU-M and Jacobi-GQU-M algorithms. In this subsection, we propose two natural variants of Jacobi-GLU and Jacobi-GQU algorithms, which will be called Jacobi-GLU-M and Jacobi-GQU-M algorithms, respectively. In these two algorithms, in each iteration, among all the index pairs $\left(i_{k}, j_{k}\right)$ and elementary functions $\nu_{k}$ satisfying (57), we choose $\left(i_{k}, j_{k}\right)$ and $\nu_{k}$ such that the cost function obtains the largest reduction.
7.3.3. Global convergence. Similar as in Section 7.1, we have the following result about the global convergence of Jacobi-GLU and Jacobi-GQU algorithms for cost function (6). It is clear that Theorem 7.5 also applies to Jacobi-GLU-M and Jacobi-GQU-M algorithms.

Theorem 7.5. In Algorithm 2 and Algorithm 3 for cost function (6), if condition (18) is always satisfied, then the iterates $\left\{\boldsymbol{X}_{k}\right\}_{k \geq 1}$ converge to a point $\boldsymbol{X}_{*}$.

[^9]```
Algorithm 2: Jacobi-GLU
    Input: A starting point \(\boldsymbol{X}_{0}\), a positive constant \(0<\varepsilon<\sqrt{\frac{2}{3 m(m-1)}}\).
    Output: Sequence of iterates \(\left\{\boldsymbol{X}_{k}\right\}_{k \geq 1}\).
    for \(k=1,2, \cdots\), do
        Choose the index pair \(\left(i_{k}, j_{k}\right)\) and \(\nu_{k}\) such that
\[
\begin{equation*}
\left\|\partial \nu_{k}\left(\boldsymbol{I}_{2}\right)\right\| \geq \varepsilon\left\|\boldsymbol{\Lambda}\left(\boldsymbol{X}_{k-1}\right)\right\| \tag{57}
\end{equation*}
\]
        where \(\nu_{k}=\varphi_{k}, \psi_{k}\) or \(\rho_{k}\);
        Compute \(\boldsymbol{\Psi}_{k}^{*}\) that minimizes the function \(\nu_{k}\);
        Update \(\boldsymbol{X}_{k}=\boldsymbol{X}_{k-1} \boldsymbol{V}_{k}\), where \(\boldsymbol{V}_{k}=\boldsymbol{U}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}, \boldsymbol{L}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}\) or \(\boldsymbol{D}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}\).
    end for
```

```
Algorithm 3: Jacobi-GQU
```

Algorithm 3: Jacobi-GQU
Input: A starting point $\boldsymbol{X}_{0}$, a positive constant $0<\varepsilon<\sqrt{\frac{3+\sqrt{5}}{3 m(m-1)}}$.
Input: A starting point $\boldsymbol{X}_{0}$, a positive constant $0<\varepsilon<\sqrt{\frac{3+\sqrt{5}}{3 m(m-1)}}$.
Output: Sequence of iterates $\left\{\boldsymbol{X}_{k}\right\}_{k \geq 1}$.
Output: Sequence of iterates $\left\{\boldsymbol{X}_{k}\right\}_{k \geq 1}$.
for $k=1,2, \cdots$, do
for $k=1,2, \cdots$, do
Choose the index pair $\left(i_{k}, j_{k}\right)$ and $\nu_{k}$ satisfying (57), where $\nu_{k}=h_{k}, \psi_{k}$ or $\rho_{k}$;
Choose the index pair $\left(i_{k}, j_{k}\right)$ and $\nu_{k}$ satisfying (57), where $\nu_{k}=h_{k}, \psi_{k}$ or $\rho_{k}$;
Compute $\boldsymbol{\Psi}_{k}^{*}$ that minimizes the function $\nu_{k}$;
Compute $\boldsymbol{\Psi}_{k}^{*}$ that minimizes the function $\nu_{k}$;
Update $\boldsymbol{X}_{k}=\boldsymbol{X}_{k-1} \boldsymbol{V}_{k}$, where $\boldsymbol{V}_{k}=\boldsymbol{G}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}, \boldsymbol{L}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}$ or $\boldsymbol{D}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}$.
Update $\boldsymbol{X}_{k}=\boldsymbol{X}_{k-1} \boldsymbol{V}_{k}$, where $\boldsymbol{V}_{k}=\boldsymbol{G}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}, \boldsymbol{L}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}$ or $\boldsymbol{D}^{\left(i_{k}, j_{k}, \boldsymbol{\Psi}_{k}^{*}\right)}$.
end for

```
    end for
```

7.3.4. Weak convergence. We also have the following result about the weak convergence of Jacobi-GLU and Jacobi-GQU algorithms for cost function (6).

Theorem 7.6. In Algorithm 2 and Algorithm 3 for cost function (6), if condition (18) is always satisfied, and $\boldsymbol{X}_{*}$ is an accumulation point of the iterates $\left\{\boldsymbol{X}_{k}\right\}_{k \geq 1}$, then $\boldsymbol{X}_{*}$ is a stationary point of the cost function (6).

Remark 7.7. In Algorithm 2 and Algorithm 3, a more natural way of choosing the index pair $\left(i_{k}, j_{k}\right)$ is according to a cyclic ordering. In fact, this cyclic way has been often used in the literature $[35,41,45]$. In this case, we call them the Jacobi-CLU (based on LU decomposition) and Jacobi-CQU (based on QU decomposition) algorithms, respectively. To our knowledge, there has been no theoretical result about the convergence of Jacobi-CLU and Jacobi-CQU algorithms in the literature. In this paper, we propose Jacobi-GLU and Jacobi-GQU algorithms, and establish their convergence. It may be also interesting to study how to establish the convergence of Jacobi-CLU and Jacobi-CQU algorithms.

## 8. Experiments

In this section, we conduct some numerical experiments to compare the performances of Jacobi-GLU in Algorithm 2, Jacobi-GLU-M in Section 7.3.2, Jacobi-CLU in Remark 7.7, JacobiGQU in Algorithm 3, Jacobi-GQU-M in Section 7.3.2, Jacobi-CQU in Remark 7.7, Jacobi-GQ and Jacobi-CQ algorithms. Here, we denote by Jacobi-GQ the gradient based Jacobi-type
algorithm on unitary group proposed in [44], and by Jacobi-CQ the Jacobi-type algorithm on unitary group with a cyclic ordering. Note that Jacobi-GQ and Jacobi-CQ find the iterations only in $\mathbf{U}_{n}$, not in $\mathbf{S L}_{n}(\mathbb{C})$. All the algorithms stop after 1000 iterations.

Example 8.1. For the following sets of complex matrices, we set $\mu_{\ell}=1$ and $(\cdot)^{\boldsymbol{\wedge}}=(\cdot)^{\mathrm{H}}$ in the cost function (6), and run the eight Jacobi-type algorithms to minimize it. The values of cost function (6) in the iterations are shown in Figure 1.
(i) We randomly generate two complex matrices $\left\{\boldsymbol{A}_{\ell}\right\}_{1 \leq \ell \leq 2} \subseteq \mathbb{C}^{5 \times 5}$.
(ii) We randomly generate a complex matrix $\boldsymbol{X} \in \mathbb{C}^{10 \times 10}$, and set $\boldsymbol{A}^{(\ell)}=\boldsymbol{X}^{\mathrm{H}}\left(\boldsymbol{I}_{10}+\boldsymbol{e}_{\ell}^{\top} \boldsymbol{e}_{\ell}\right) \boldsymbol{X}$ for $1 \leq \ell \leq 10$.
(iii) We randomly generate a complex upper triangular matrix $\boldsymbol{X} \in \mathbf{U T}_{10}(\mathbb{C})$, complex diagonal matrices $\left\{\boldsymbol{D}_{\ell}\right\}_{1 \leq \ell \leq 10} \subseteq \mathbb{C}^{10 \times 10}$, and set $\boldsymbol{A}^{(\ell)}=\boldsymbol{X}^{\mathrm{H}} \boldsymbol{D}_{\ell} \boldsymbol{X}$ for $1 \leq \ell \leq 10$.
(iv) We randomly generate a complex nonsingular matrix $\boldsymbol{X} \in \mathbf{S L}_{10}(\mathbb{C})$, complex diagonal matrices $\left\{\boldsymbol{D}_{\ell}\right\}_{1 \leq \ell \leq 10} \subseteq \mathbb{C}^{10 \times 10}$, and set $\boldsymbol{A}^{(\ell)}=\boldsymbol{X}^{\mathrm{H}} \boldsymbol{D}_{\ell} \boldsymbol{X}$ for $1 \leq \ell \leq 10$.

Example 8.2. Let $\left\{\boldsymbol{A}_{\ell}\right\}_{1 \leq \ell \leq L} \subseteq \mathbb{C}^{n \times n}$ be randomly generated with different $n$ and $L$ values as shown in Table 2. Let $\mu_{\ell}=1$ and $(\cdot)^{\wedge}=(\cdot)^{\mathrm{H}}$ in the cost function (6). We run the eight Jacobi-type algorithms to minimize it, and get the final values of cost function (6) in Table 2. In this table, the numbers appearing in bold indicate that this number is the best result in the row where it lies.

TABLE 2. Cost function values in Example 8.2

|  |  | GLU-M | GLU | CLU | GQU-M | GQU | CQU | GQ | CQ |
| :--- | :--- | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $n=5$ | $L=1$ | 0.0670 | 0.0133 | $\mathbf{0 . 0 0 0 0}$ | 0.0229 | 0.0921 | 0.0692 | 0.8098 | 0.8098 |
| $n=5$ | $L=2$ | $\mathbf{1 . 1 9 3 3}$ | 1.8995 | 1.9069 | 2.0520 | 2.1803 | 2.3514 | 3.3850 | 3.3850 |
| $n=5$ | $L=5$ | 11.3882 | 11.3908 | 11.3967 | $\mathbf{1 1 . 3 8 8 0}$ | 11.3883 | 11.4016 | 13.3835 | 13.3835 |
| $n=10$ | $L=1$ | 0.3070 | $\mathbf{0 . 2 2 2 0}$ | 0.3229 | 0.6100 | 0.6640 | 1.5259 | 3.2546 | 3.2546 |
| $n=10$ | $L=5$ | $\mathbf{4 7 . 5 9 9 3}$ | 48.5667 | 50.2540 | 47.6150 | 48.3432 | 48.1588 | 54.0629 | 53.8712 |
| $n=10$ | $L=10$ | 119.0841 | 120.8933 | 121.9417 | $\mathbf{1 1 5 . 7 7 4 8}$ | 117.8872 | 119.3527 | 123.2703 | 123.2546 |

From these experimental results, we can see that, compared with the Jacobi-GQ and JacobiCQ algorithms on unitary group, the Jacobi-type algorithms on $\mathbf{S L}_{m}(\mathbb{C})$ considered in this paper always obtain smaller cost function values, and they need more iterations to attain the best cost functions values. Moreover, compared with the QU decomposition based Jacobi-type algorithms, the LU decomposition based ones obtain better experimental results in most cases.

## 9. Conclusions

In this paper, to solve JADM problem (3), which is important in blind source separation, we formulate two different equivalent formulations, i.e. problem (5) defined on $\mathbf{S t}(m, n, \mathbb{C}) \times$ $\mathbf{S L}_{m}(\mathbb{C})$, and problem (6) defined on $\mathbf{S L}_{m}(\mathbb{C})$. Then, for these two approaches, we propose three


Figure 1. Cost function values in Example 8.1.

BCD-G algorithms and two Jacobi-type algorithms, and establish their global and weak convergence. In the future work, it may be interesting to extend these algorithms and convergence results to the higher order tensor cases.

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[^1]:    ${ }^{1}$ This equation means that a complex matrix $\boldsymbol{Z} \in \mathbf{R S L}(m, n, \mathbb{C})$ if and only if there exist $\boldsymbol{Y} \in \mathbf{S t}(m, n, \mathbb{C})$ and $\boldsymbol{X} \in \mathbf{S L}_{m}(\mathbb{C})$ such that $\boldsymbol{Z}=\boldsymbol{Y} \boldsymbol{X}$. It is the same case with equations (11), (12), (13), (15) and (16).

[^2]:    ${ }^{2}$ See [3, Section 3.6] for a detailed definition.

[^3]:    ${ }^{3}$ The inequality (10) can be seen as a block coordinate analogue of [20, Eq. (3.3)] and [24, Eq. (10)].
    ${ }^{4}$ The reason why we use plane diagonal transformations will be shown in Section 4

[^4]:    ${ }^{5}$ For any starting point, the iterates converge as a whole sequence.
    ${ }^{6}$ Every accumulation point is a stationary point.

[^5]:    ${ }^{7}$ See [21, Definition 2.7.1] or [25, Definition 5.1] for a definition of an analytic submanifold.
    ${ }^{8}$ The values of $\delta, \sigma, \zeta$ depend on the specific point in question.

[^6]:    ${ }^{9}$ It is also known as the first Wolfe condition in the literature.

[^7]:    ${ }^{10}$ The inequality (38) can be seen as a non-orthogonal analogue of [20, Eq. (3.3)] and [24, Eq. (10)].

[^8]:    ${ }^{11}$ In the $(\cdot)^{\star}=(\cdot)^{\mathrm{H}}$ case, this expression was first formulated in [12].

[^9]:    ${ }^{12}$ In this case, we do not consider the GU class of transformations.
    ${ }^{13}$ These algorithms are based on the similar ideas as the Jacobi-G type algorithms in [20, 24, 44].

