# AN EFFICIENT ITERATIVE METHOD FOR SOLVING QUADRATIC INVERSE EIGENVALUE PROBLEM OVER HERMITIAN CENTROSKEW MATRICES WITH A SUBMATRIX CONSTRAINT* 

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#### Abstract

In this paper, a nonlinear conjugate gradient method is constructed to find the least squares solution of the quadratic inverse eigenvalue problem under a prescribed submatrix constraint. The convergence analysis of the proposed method is given. A method for choosing the initial matrix group to obtain the least Frobenius norm solution is provided and the unique optimal approximation solution to a given matrix group is derived. Numerical examples are provided to illustrate the efficiency of the proposed method and testify the conclusions suggested in this paper.


Key words: quadratic inverse eigenvalue problem, Hermitian centroskew matrix, submatrix constraint
Mathematics Subject Classification: 15A24, 15A57, 65F10, 65F30

## 1 Introduction

Throughout this paper, the following notations and concepts are needed. Let $\mathbb{C}^{m \times n}$ be the set of all complex $m \times n$ matrices. For $A \in \mathbb{C}^{m \times n}$, we write $A^{T}, A^{H}, \operatorname{tr}(A)$ and $\mathcal{R}(A)$ to denote the transpose, the conjugate transpose, the trace and the column space of matrix $A$, respectively. Let $I_{n}=\left(e_{1}, e_{2}, \ldots, e_{n}\right)$ and $S_{n}=\left(e_{n}, e_{n-1}, \ldots, e_{1}\right)$ be the $n \times n$ identity matrix and reverse identity matrix, respectively, where $e_{i}$ denotes the $i$ th column of the identity matrix $I_{n}$. In the space $\mathbb{C}^{m \times n}$ over the real number field $\mathbb{R}$, the inner product for $A, B$ is $\langle A, B\rangle=\operatorname{Re}\left[\operatorname{tr}\left(B^{H} A\right)\right]$. Then $\|A\|^{2}=\langle A, A\rangle=\operatorname{Re}\left[\operatorname{tr}\left(A^{H} A\right)\right]=\operatorname{tr}\left(A^{H} A\right)=\|A\|_{F}^{2}$ and $\|A \pm B\|^{2}=\|A\|^{2}+\|B\|^{2} \pm 2\langle A, B\rangle$ for $A, B \in \mathbb{C}^{m \times n}$. We say $A$ and $B$ are orthogonal if $\langle A, B\rangle=0$. Let $D_{p, n}=\left\{d=\left(d_{1}, d_{2}, \ldots, d_{p}\right): 1 \leq d_{1}<d_{2}<\cdots<d_{p} \leq n\right\}$ be the strictly increasing sequences of $p$ elements from $1,2, \ldots, n$. For $s=\left(s_{1}, s_{2}, \ldots, s_{p}\right) \in D_{p, n}, t=$ $\left(t_{1}, t_{2}, \ldots, t_{q}\right) \in D_{q, n}$ and $u=\left(u_{1}, u_{2}, \ldots, u_{r}\right) \in D_{r, n}$, we denote $E_{s}=\left(e_{s_{1}}, e_{s_{2}}, \ldots, e_{s_{p}}\right) \in$ $\mathbb{C}^{n \times p}, E_{t}=\left(e_{t_{1}}, e_{t_{2}}, \ldots, e_{t_{q}}\right) \in \mathbb{C}^{n \times q}$ and $E_{u}=\left(e_{u_{1}}, e_{u_{2}}, \ldots, e_{u_{r}}\right) \in \mathbb{C}^{n \times r}$. Let $A[s \mid t]$ be the $p \times q$ submatrix of $A \in \mathbb{C}^{m \times n}$ determined by rows indexed by $s$ and columns indexed by $t$, and let $A[\bar{s} \mid \bar{t}]$ be the $(m-p) \times(n-q)$ submatrix of $A \in \mathbb{C}^{m \times n}$ determined by deleting rows indexed by $s$ and columns indexed by $t$. A matrix $A \in \mathbb{C}^{n \times n}$ is said to be a Hermitian centroskew matrix if $\bar{a}_{i j}=a_{j i}$ and $a_{i j}=-a_{n+1-i, n+1-j}$ for all $1 \leq i, j \leq n$. Let $H S C^{n \times n}$

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be the set of all $n \times n$ Hermitian centroskew matrices. It is well known that Hermitian centroskew matrices include symmetric centroskew-symmetric and symmetric Toeplitz matrices. Hermitian centroskew matrices have practical applications in information theory, linear system theory, linear estimate theory and numerical analysis [11, 12].

In this paper, we are interested in the quadratic inverse eigenvalue problem which consists in finding matrices $A, B, C \in \mathbb{C}^{n \times n}$ such that

$$
\begin{equation*}
A X \Lambda^{2}+B X \Lambda+C X=O \tag{1.1}
\end{equation*}
$$

where $X=\left(x_{1}, x_{2}, \ldots, x_{m}\right) \in \mathbb{C}^{n \times m}$ and $\Lambda=\operatorname{diag}\left(\lambda_{1}, \lambda_{2}, \ldots, \lambda_{m}\right) \in \mathbb{C}^{m \times m}$ are given matrices. The quadratic inverse eigenvalue problem appears in some practical applications [16]. Such as, setting the pole assignment problem can be seen as a quadratic inverse eigenvalue problem [14]. The behavior of a time invariant second order differential system

$$
\begin{equation*}
A \ddot{x}(t)+B \dot{x}(t)+C x(t)=f(t) \tag{1.2}
\end{equation*}
$$

can be described by solving a quadratic inverse eigenvalue problem [35].

- If $A=O$, then the quadratic inverse eigenvalue problem (1.1) reduces to a generalized inverse eigenvalue problem

$$
\begin{equation*}
B X \Lambda+C X=O \tag{1.3}
\end{equation*}
$$

- If $A=O$ and $B=I$, then the quadratic inverse eigenvalue problem (1.1) reduces to a standard inverse eigenvalue problem

$$
\begin{equation*}
X \Lambda+C X=O \tag{1.4}
\end{equation*}
$$

In general, the inverse eigenvalue problem is as important as the eigenvalue problem $[16,30]$. The inverse eigenvalue problem has wide applications in engineering and scientific computation [39, 17, 19]. For instance, the inverse eigenvalue problem of centro-symmetric matrices with a submatrix constraint initially occurs in the design of Hopfield neural network, civil engineering and aviation [13, 3]. The parameterized generalized inverse eigenvalue problem (PGIEP) [17] arises in the discrete analogue of inverse Sturm-Liouville problem [27, 28], factor analysis [28] and structural design [22, 29, 36, 38]. The problem of designing the truss structure with specified natural frequency may be recast as a PGIEP.

The literature on solvability condition and numerical method for the inverse eigenvalue problem is large and still growing rapidly [19, 40, 9, 10, 17, 20, 21]. For instance, Bai [4] and Zhang et al. [41] derived some necessary and sufficient conditions for the inverse eigenvalue problem with Hermitian generalized skew-Hamiltonian matrices and Hermitian generalized Hamiltonian matrices, respectively. Ghanbari [19] studied the explicit expression solution of a generalized inverse eigenvalue problem, in which $C$ is a semi-infinite Jacobi matrix with positive off-diagonal entries and $B$ is a nonsingular diagonal matrix. Liu et al. [31] considered an existence condition of a solution and an analytic expression solution of the generalized inverse eigenvalue problem for centrohermitian matrices. Moghaddam et al. [33] proposed an algorithm for reconstructing penta-diagonal coefficient matrices of a generalized inverse eigenvalue problem. Wei et al. [37] discussed a generalized inverse eigenvalue problem with Hermitian generalized Hamiltonian matrices and derived an analytic expression solution by the matrix decomposition theory and Hilbert space approximation theory. Aishima [1] introduced a quadratically convergent algorithm for the inverse symmetric eigenvalue problem. Dai and Liang [18] considered the generalized inverse eigenvalue problem for
the $(P, Q)$-conjugate matrix and the associated approximation problem by using the generalized singular value decomposition and canonical correlation decomposition. Recently, Cai and Chen [9] proposed an iterative method for solving the generalized inverse eigenvalue problem and its optimal approximation problem over partially bisymmetric matrices. They also studied [10] the least squares solution of a generalized inverse eigenvalue problem over Hermitian Hamiltonian matrices with a submatrix constraint.

Inspired by the work mentioned above, we focus on the following problems.
Problem 1.1. Given $X \in \mathbb{C}^{n \times m}, \Lambda=\operatorname{diag}\left(\lambda_{1}, \lambda_{2}, \ldots, \lambda_{m}\right) \in \mathbb{C}^{m \times m}, s=\left(s_{1}, s_{2}, \ldots, s_{p}, n+\right.$ $\left.1-s_{p}, \ldots, n+1-s_{2}, n+1-s_{1}\right) \in D_{2 p, n}, t=\left(t_{1}, t_{2}, \ldots, t_{q}, n+1-t_{q}, \ldots, n+1-t_{2}, n+1-t_{1}\right) \in$ $D_{2 q, n}, u=\left(u_{1}, u_{2}, \ldots, u_{r}, n+1-u_{r}, \ldots, n+1-u_{2}, n+1-u_{1}\right) \in D_{2 r, n}, A_{p} \in \mathbb{C}^{2 p \times 2 p}$, $B_{q} \in \mathbb{C}^{2 q \times 2 q}$ and $C_{r} \in \mathbb{C}^{2 r \times 2 r}$. Let $S_{1}=\left\{Y \mid Y[s \mid s]=A_{p}, Y[\bar{s} \mid \bar{s}] \in H S C^{(n-2 p) \times(n-2 p)}\right\}$, $S_{2}=\left\{Y \mid Y[t \mid t]=B_{q}, Y[\bar{t} \mid \bar{t}] \in H S C^{(n-2 q) \times(n-2 q)}\right\}$ and $S_{3}=\left\{Y \mid Y[u \mid u]=C_{r}, Y[\bar{u} \mid \bar{u}] \in\right.$ $\left.H S C^{(n-2 r) \times(n-2 r)}\right\}$. We want to find $A_{*} \in S_{1}, B_{*} \in S_{2}$ and $C_{*} \in S_{3}$ such that

$$
\left\|A_{*} X \Lambda^{2}+B_{*} X \Lambda+C_{*} X\right\|=\min _{(A, B, C) \in S_{1} \times S_{2} \times S_{3}}\left\|A X \Lambda^{2}+B X \Lambda+C X\right\|
$$

Problem 1.2. Let $S_{E}$ be the set of solutions of Problem 1.1. For given $\bar{A}, \bar{B}, \bar{C} \in \mathbb{C}^{n \times n}$, we want to find $(\widehat{A}, \widehat{B}, \widehat{C}) \in S_{E}$ such that

$$
\|\widehat{A}-\bar{A}\|^{2}+\|\widehat{B}-\bar{B}\|^{2}+\|\widehat{C}-\bar{C}\|^{2}=\min _{(A, B, C) \in S_{E}}\left[\|A-\bar{A}\|^{2}+\|B-\bar{B}\|^{2}+\|C-\bar{C}\|^{2}\right]
$$

Although Hajarian [23, 25] constructed the conjugate direction and BCR methods to solve some special quadratic inverse eigenvalue problems. Hajarian and Hassan [24] established the CGNR method for finding the least squares solution of a quadratic inverse eigenvalue problem with partially bisymmetric matrices under a prescribed submatrix constraint. In fact, Hajarian [23, 25] and Hajarian and Hassan [24] extended the subspace method for linear system to the setting of quadratic inverse eigenvalue problem from the perspective of numerical algebra. Naturally, we want to know whether the quadratic inverse eigenvalue problem can be solved from the perspective of optimization or not. This is one of our motivations in this paper.

Since matrices $X$ and $\Lambda$ are known in the system (1.1), the investigated problem (Problem 1.1) can be transformed into finding the least squares solution of the generalized Sylvester matrix equation (1.1). Such a problem has been investigated in many papers. For example, Beik and Salkuyeh [5, 6] established some modified conjugate gradient type methods. However, the proposed methods in $[5,6]$ are extensions of the classical subspace methods for linear system and feasible only when the considered generalized Sylvester matrix equation is consistent. They can not be directly applied for the matrix equation with a submatrix constraint. Naturally, one may ask :"does the nonlinear conjugate gradient method can be extended to solve the quadratic inverse eigenvalue problem?" As another motivation in this paper, we are interested in studying a nonlinear conjugate gradient method for solving the quadratic inverse eigenvalue problem with a submatrix constraint. The contributions of this paper are as follows.

- We establish a nonlinear conjugate gradient method for finding the least squares solution of the quadratic inverse eigenvalue problem.
- We give the convergence analysis of the proposed nonlinear conjugate gradient method. A method for choosing the initial matrices to obtain the least Frobenius norm least squares solution of the quadratic inverse eigenvalue problem is given.
- By reformulating Problem 1.2 as another quadratic inverse eigenvalue problem with a submatrix constraint, we apply the nonlinear conjugate gradient method to solve it.
- Experimental results demonstrate that our proposed method yields superior performance for the quadratic inverse eigenvalue problem.

The remainder of this paper is organized as follows. In Section 2, we give an equivalent characterization of Problem 1.1 and establish a nonlinear conjugate gradient method to solve it. In Section 3, we present an iterative method for solving Problem 1.2. In Section 4, we report some numerical results to illustrate the feasibility and efficiency of the proposed method. This paper ends up with some concluding remarks in Section 5.

## 2 Iterative Method for Problem 1.1

Notice that the solution set of Problem 1.1 is merely a linear manifold rather than a linear subspace. First, we give an equivalent characterization of Problem 1.1. Let

$$
\begin{aligned}
S_{1} & =\left\{Y \mid Y[s \mid s]=A_{p}, Y[\bar{s} \mid \bar{s}]=O\right\} \oplus \widetilde{S}_{1} \\
S_{2} & =\left\{Y \mid Y[t \mid t]=B_{q}, Y[\bar{t} \mid \bar{t}]=O\right\} \oplus \widetilde{S}_{2} \\
S_{3} & =\left\{Y \mid Y[u \mid u]=C_{r}, Y[\bar{u} \mid \bar{u}]=O\right\} \oplus \widetilde{S}_{3}
\end{aligned}
$$

where

$$
\begin{align*}
\widetilde{S}_{1} & =\left\{Y \mid Y \in H S C^{n \times n}, Y[s \mid s]=O\right\}  \tag{2.1a}\\
\widetilde{S}_{2} & =\left\{Y \mid Y \in H S C^{n \times n}, Y[t \mid t]=O\right\}  \tag{2.1b}\\
\widetilde{S}_{3} & =\left\{Y \mid Y \in H S C^{n \times n}, Y[u \mid u]=O\right\} . \tag{2.1c}
\end{align*}
$$

Then Problem 1.1 has the following equivalent form.
Problem 2.1. Given $X \in \mathbb{C}^{n \times m}, \Lambda=\operatorname{diag}\left(\lambda_{1}, \lambda_{2}, \ldots, \lambda_{m}\right) \in \mathbb{C}^{m \times m}, s=\left(s_{1}, s_{2}, \ldots, s_{p}, n+\right.$ $\left.1-s_{p}, \ldots, n+1-s_{2}, n+1-s_{1}\right) \in D_{2 p, n}, t=\left(t_{1}, t_{2}, \ldots, t_{q}, n+1-t_{q}, \ldots, n+1-t_{2}, n+1-t_{1}\right) \in$ $D_{2 q, n}, u=\left(u_{1}, u_{2}, \ldots, u_{r}, n+1-u_{r}, \ldots, n+1-u_{2}, n+1-u_{1}\right) \in D_{2 r, n}, A_{p} \in \mathbb{C}^{2 p \times 2 p}$, $B_{q} \in \mathbb{C}^{2 q \times 2 q}$ and $C_{r} \in \mathbb{C}^{2 r \times 2 r}$. We want to find $\widetilde{A}_{*} \in \widetilde{S}_{1}, \widetilde{B}_{*} \in \widetilde{S}_{2}$ and $\widetilde{C}_{*} \in \widetilde{S}_{3}$ such that

$$
\left\|\widetilde{A}_{*} X \Lambda^{2}+\widetilde{B}_{*} X \Lambda+\widetilde{C}_{*} X-\widetilde{Z}\right\|=\min _{(\widetilde{A}, \widetilde{B}, \widetilde{C}) \in \widetilde{S}_{1} \times \widetilde{S}_{2} \times \widetilde{S}_{3}}\left\|\widetilde{A} X \Lambda^{2}+\widetilde{B} X \Lambda+\widetilde{C} X-\widetilde{Z}\right\|
$$

where $\widetilde{Z}=-\widetilde{A}_{p} X \Lambda^{2}-\widetilde{B}_{q} X \Lambda-\widetilde{C}_{r} X$, in which $\widetilde{A}_{p}, \widetilde{B}_{q}$ and $\widetilde{C}_{r}$ denote the matrices satisfying $\widetilde{A}_{p}[s \mid s]=A_{p}, \widetilde{B}_{q}[t \mid t]=B_{q}, \widetilde{C}_{r}[u \mid u]=C_{r}$ and zeros elsewhere.

It is easy to see that $\left(\widetilde{A}_{*}, \widetilde{B}_{*}, \widetilde{C}_{*}\right)$ is a solution of Problem 2.1 if and only if $(A, B, C)=$ $\left(\widetilde{A}_{*}+\widetilde{A}_{p}, \widetilde{B}_{*}+\widetilde{B}_{q}, \widetilde{C}_{*}+\widetilde{C}_{r}\right)$ is a solution of Problem 1.1.

In the sequel, we illustrate that the solution to Problem 2.1 exists. First, the property of Hermitian centroskew matrix is needed.

Lemma 2.1 ([27]). A matrix $X \in H S C^{n \times n}$ if and only if $X=X^{H}=-S_{n} X S_{n}$.
The following result is important for the proof of the existence of the solution to Problem 2.1.

Lemma 2.2 ([26]). If a quadratic function $f$ has a lower bound in a nonempty polyhedron $\mathbb{X}$, then $f$ must have a global minimum in $\mathbb{X}$.

Proposition 2.3. Let $L_{1}(\widetilde{A}, \widetilde{B}, \widetilde{C})=\widetilde{A} X \Lambda^{2}+\widetilde{B} X \Lambda+\widetilde{C} X-\widetilde{Z}, L_{2}(\widetilde{A}, \widetilde{B}, \widetilde{C})=\widetilde{A}^{H} X \Lambda^{2}+$ $\widetilde{B}^{H} X \Lambda+\widetilde{C}^{H} X-\widetilde{Z}, L_{3}(\widetilde{A}, \widetilde{B}, \widetilde{C})=S_{n} \widetilde{A} S_{n} X \Lambda^{2}+S_{n} \widetilde{B} S_{n} X \Lambda+S_{n} \widetilde{C} S_{n} X+\widetilde{Z}$ and $L_{4}(\widetilde{A}, \widetilde{B}, \widetilde{C})=$ $S_{n} \widetilde{A}^{H} S_{n} X \Lambda^{2}+S_{n} \widetilde{B}^{H} S_{n} X \Lambda+S_{n} \widetilde{C}^{H} S_{n} X+\widetilde{Z}$. Then

$$
\begin{align*}
& \left\{\widetilde{A} \in \mathbb{C}^{n \times n}, \widetilde{B} \in \mathbb{C}^{n \times n}, \widetilde{C} \in \mathbb{C}^{n \times n}: \widetilde{A}[s \mid s]=O, \widetilde{B}[t \mid t]=O, \widetilde{C}[u \mid u]=O\right\} \\
&  \tag{2.2}\\
& \quad\left\|\frac{L_{1}(\widetilde{A}, \widetilde{B}, \widetilde{C})}{4}+\frac{L_{2}(\widetilde{A}, \widetilde{B}, \widetilde{C})}{4}-\frac{L_{3}(\widetilde{A}, \widetilde{B}, \widetilde{C})}{4}-\frac{L_{4}(\widetilde{A}, \widetilde{B}, \widetilde{C})}{4}\right\|^{2}
\end{align*}
$$

takes a global minimum value at $\left(\widetilde{A}_{*}, \widetilde{B}_{*}, \widetilde{C}_{*}\right)$, where $\widetilde{A}_{*}[s \mid s]=O, \widetilde{B}_{*}[t \mid t]=O$ and $\widetilde{C}_{*}[u \mid u]=$ $O$.

Proof. Let

$$
\begin{aligned}
& h(\operatorname{vec}(\widetilde{A}), \operatorname{vec}(\widetilde{B}), \operatorname{vec}(\widetilde{C})) \\
= & \left\|\frac{1}{4}\left\{\operatorname{vec}\left[L_{1}(\widetilde{A}, \widetilde{B}, \widetilde{C})\right]+\operatorname{vec}\left[L_{2}(\widetilde{A}, \widetilde{B}, \widetilde{C})\right]-\operatorname{vec}\left[L_{3}(\widetilde{A}, \widetilde{B}, \widetilde{C})\right]-\operatorname{vec}\left[L_{4}(\widetilde{A}, \widetilde{B}, \widetilde{C})\right]\right\}\right\|^{2}
\end{aligned}
$$

Then $h$ is a quadratic function about variable $(\operatorname{vec}(\widetilde{A}), \operatorname{vec}(\widetilde{B}), \operatorname{vec}(\widetilde{C}))$. Let

$$
\mathbb{X}=\left\{\operatorname{vec}(\widetilde{A}) \in \mathbb{C}^{n^{2}}, \operatorname{vec}(\widetilde{B}) \in \mathbb{C}^{n^{2}}, \operatorname{vec}(\widetilde{C}) \in \mathbb{C}^{n^{2}}: \widetilde{A}[s \mid s]=O, \widetilde{B}[t \mid t]=O, \widetilde{C}[u \mid u]=O\right\}
$$

Obviously, $\mathbb{X}$ is a nonempty polyhedron. It then follows from Lemma 2.2 that the problem (2.2) takes the global minimum value at some point $\left(\operatorname{vec}\left(\widetilde{A}_{*}\right), \operatorname{vec}\left(\widetilde{B}_{*}\right), \operatorname{vec}\left(\widetilde{C}_{*}\right)\right)$, where $\widetilde{A}_{*}[s \mid s]=O, \widetilde{B}_{*}[t \mid t]=O$ and $\widetilde{C}_{*}[u \mid u]=O$. By the inverse vec operator, the result is established. The proof is completed.

Proposition 2.4. The solution to Problem 2.1 exists.
Proof. For arbitrary $\widetilde{A}, \widetilde{B}, \widetilde{C} \in H S C^{n \times n}$, it follows from Lemma 2.1 that

$$
\begin{aligned}
& \left\|\widetilde{A} X \Lambda^{2}+\widetilde{B} X \Lambda+\widetilde{C} X-\widetilde{Z}\right\|^{2} \\
& = \\
& =\| \frac{\widetilde{A}+\widetilde{A}^{H}-S_{n} \widetilde{A} S_{n}-S_{n} \widetilde{A}^{H} S_{n}}{4} X \Lambda^{2} \\
& \quad+\frac{\widetilde{B}+\widetilde{B}^{H}-S_{n} \widetilde{B} S_{n}-S_{n} \widetilde{B}^{H} S_{n}}{4} X \Lambda+\frac{\widetilde{C}+\widetilde{C}^{H}-S_{n} \widetilde{C} S_{n}-S_{n} \widetilde{C}^{H} S_{n}}{4} X-\widetilde{Z} \|^{2} \\
& =
\end{aligned} \quad\left\|\frac{L_{1}(\widetilde{A}, \widetilde{B}, \widetilde{C})}{4}+\frac{L_{2}(\widetilde{A}, \widetilde{B}, \widetilde{C})}{4}-\frac{L_{3}(\widetilde{A}, \widetilde{B}, \widetilde{C})}{4}-\frac{L_{4}(\widetilde{A}, \widetilde{B}, \widetilde{C})}{4}\right\|^{2} .
$$

Assume that the problem (2.2) takes the global minimum value at $\left(\widetilde{A}_{*}, \widetilde{B}_{*}, \widetilde{C}_{*}\right)$, where $\widetilde{A}_{*}[s \mid s]=O, \widetilde{B}_{*}[t \mid t]=O$ and $\widetilde{C}_{*}[u \mid u]=O$. Then, for arbitrary $(\widetilde{A}, \widetilde{B}, \widetilde{C}) \in \widetilde{S}_{1} \times \widetilde{S}_{2} \times \widetilde{S}_{3}$, we
have

$$
\begin{align*}
& \left\|\widetilde{A} X \Lambda^{2}+\widetilde{B} X \Lambda+\widetilde{C} X-\widetilde{Z}\right\|^{2} \\
& =\left\|\frac{L_{1}(\widetilde{A}, \widetilde{B}, \widetilde{C})}{4}+\frac{L_{2}(\widetilde{A}, \widetilde{B}, \widetilde{C})}{4}-\frac{L_{3}(\widetilde{A}, \widetilde{B}, \widetilde{C})}{4}-\frac{L_{4}(\widetilde{A}, \widetilde{B}, \widetilde{C})}{4}\right\|^{2} \\
& \geq \min _{\left\{\widetilde{A} \in \mathbb{C}^{n \times n}, \widetilde{B} \in \mathbb{C}^{n \times n}, \widetilde{C} \in \mathbb{C}^{n \times n}: \widetilde{A}[s \mid s]=O, \widetilde{B}[t \mid t]=O, \widetilde{C}[u \mid u]=O\right\}} \| \frac{L_{1}(\widetilde{A}, \widetilde{B}, \widetilde{C})}{4}+\frac{L_{2}(\widetilde{A}, \widetilde{B}, \widetilde{C})}{4} \\
& \quad-\frac{L_{3}(\widetilde{A}, \widetilde{B}, \widetilde{C})}{4}-\frac{L_{4}(\widetilde{A}, \widetilde{B}, \widetilde{C})}{4} \|^{2} \\
& =\left\|\frac{L_{1}\left(\widetilde{A}_{*}, \widetilde{B}_{*}, \widetilde{C}_{*}\right)}{4}+\frac{L_{2}\left(\widetilde{A}_{*}, \widetilde{B}_{*}, \widetilde{C}_{*}\right)}{4}-\frac{L_{3}\left(\widetilde{A}_{*}, \widetilde{B}_{*}, \widetilde{C}_{*}\right)}{4}-\frac{L_{4}\left(\widetilde{A}_{*}, \widetilde{B}_{*}, \widetilde{C}_{*}\right)}{4}\right\|^{2} \\
& =\| \frac{\widetilde{A}_{*}+\widetilde{A}_{*}^{H}-S_{n} \widetilde{A}_{*} S_{n}-S_{n} \widetilde{A}_{*}^{H} S_{n}}{4} X \Lambda^{2}+\frac{\widetilde{B}_{*}+\widetilde{B}_{*}^{H}-S_{n} \widetilde{B}_{*} S_{n}-S_{n} \widetilde{B}_{*}^{H} S_{n}}{4} X \Lambda \\
& \quad+\frac{\widetilde{C}_{*}+\widetilde{C}_{*}^{H}-S_{n} \widetilde{C}_{*} S_{n}-S_{n} \widetilde{C}_{*}^{H} S_{n}}{4} X-\widetilde{Z} \|^{2} . \tag{2.3}
\end{align*}
$$

Let

$$
\begin{aligned}
\widetilde{\mathcal{A}}_{*} & =\frac{\widetilde{A}_{*}+\widetilde{A}_{*}^{H}-S_{n} \widetilde{A}_{*} S_{n}-S_{n} \widetilde{A}_{*}^{H} S_{n}}{4} \\
\widetilde{\mathcal{B}}_{*} & =\frac{\widetilde{B}_{*}+\widetilde{B}_{*}^{H}-S_{n} \widetilde{B}_{*} S_{n}-S_{n} \widetilde{B}_{*}^{H} S_{n}}{4} \\
\widetilde{\mathcal{C}}_{*} & =\frac{\widetilde{C}_{*}+\widetilde{C}_{*}^{H}-S_{n} \widetilde{C}_{*} S_{n}-S_{n} \widetilde{C}_{*}^{H} S_{n}}{4}
\end{aligned}
$$

Using Lemma 2.1 again and the fact that $\widetilde{A}_{*}[s \mid s]=O, \widetilde{B}_{*}[t \mid t]=O$ and $\widetilde{C}_{*}[u \mid u]=O$, we obtain

$$
\widetilde{\mathcal{A}}_{*} \in \widetilde{S}_{1}, \widetilde{\mathcal{B}}_{*} \in \widetilde{S}_{2}, \widetilde{\mathcal{C}}_{*} \in \widetilde{S}_{3}
$$

It then follows from (2.3) that

$$
\min _{(\widetilde{A}, \widetilde{B}, \widetilde{C}) \in \widetilde{S}_{1} \times \widetilde{S}_{2} \times \widetilde{S}_{3}}\left\|\widetilde{A} X \Lambda^{2}+\widetilde{B} X \Lambda+\widetilde{C} X-\widetilde{Z}\right\|=\left\|\widetilde{\mathcal{A}}_{*} X \Lambda^{2}+\widetilde{\mathcal{B}}_{*} X \Lambda+\widetilde{\mathcal{C}}_{*} X-\widetilde{Z}\right\|
$$

which implies that $\left(\widetilde{\mathcal{A}}_{*}, \widetilde{\mathcal{B}}_{*}, \widetilde{\mathcal{C}}_{*}\right)$ is a solution to Problem 2.1. The proof is completed.
Remark 2.5. From Proposition 2.4, we have that the solution to Problem 1.1 exists.
We also need the following lemmas in order to derive the normal equation of Problem 2.1.

Lemma 2.6. Let $A \in \mathbb{C}^{n \times n}$ and $X \in H S C^{n \times n}$. Then

$$
\left\langle X, \frac{A-S_{n} A S_{n}}{2}\right\rangle=\left\langle X, \frac{1}{4}\left(A+A^{H}-S_{n}\left(A+A^{H}\right) S_{n}\right)\right\rangle=\langle X, A\rangle
$$

Proof. For any $A \in \mathbb{C}^{n \times n}$ and $X \in H S C^{n \times n}$, we have

$$
\begin{aligned}
\left\langle X, \frac{A-S_{n} A S_{n}}{2}\right\rangle & =\frac{1}{2}\langle X, A\rangle-\frac{1}{2}\left\langle X, S_{n} A S_{n}\right\rangle=\frac{1}{2}\langle X, A\rangle-\frac{1}{2}\left\langle S_{n} X S_{n}, A\right\rangle \\
& =\frac{1}{2}\langle X, A\rangle-\frac{1}{2}\langle-X, A\rangle=\frac{1}{2}\langle X, A\rangle+\frac{1}{2}\langle X, A\rangle=\langle X, A\rangle .
\end{aligned}
$$

Then

$$
\begin{aligned}
& \left\langle X, \frac{\left(A+A^{H}\right)-S_{n}\left(A+A^{H}\right) S_{n}}{4}\right\rangle=\frac{1}{2}\left(\left\langle X, \frac{A-S_{n} A S_{n}}{2}\right\rangle+\left\langle X, \frac{A^{H}-S_{n} A^{H} S_{n}}{2}\right\rangle\right) \\
= & \frac{1}{2}\left(\left\langle X, \frac{A-S_{n} A S_{n}}{2}\right\rangle+\left\langle X, \frac{A-S_{n} A S_{n}}{2}\right\rangle\right)=\frac{1}{2}(\langle X, A\rangle+\langle X, A\rangle)=\langle X, A\rangle .
\end{aligned}
$$

The proof is completed.
Lemma 2.7. Suppose that $X \in \mathbb{C}^{n \times n}, W_{1} \in \widetilde{S}_{1}, W_{2} \in \widetilde{S}_{2}$ and $W_{3} \in \widetilde{S}_{3}$. Then $\left\langle E_{s} E_{s}^{T} X E_{s} E_{s}^{T}, W_{1}\right\rangle=0,\left\langle E_{t} E_{t}^{T} X E_{t} E_{t}^{T}, W_{2}\right\rangle=0$ and $\left\langle E_{u} E_{u}^{T} X E_{u} E_{u}^{T}, W_{3}\right\rangle=0$.

Proof. Since $W_{1} \in \widetilde{S}_{1}, W_{2} \in \widetilde{S}_{2}$ and $W_{3} \in \widetilde{S}_{3}$, we immediately have $E_{s}^{T} W_{1} E_{s}=0$, $E_{t}^{T} W_{2} E_{t}=0$ and $E_{u}^{T} W_{3} E_{u}=0$. Hence

$$
\begin{aligned}
\left\langle E_{s} E_{s}^{T} X E_{s} E_{s}^{T}, W_{1}\right\rangle & =\operatorname{Re}\left\{\operatorname{tr}\left(W_{1}^{H} E_{s} E_{s}^{T} X E_{s} E_{s}^{T}\right)\right\} \\
& =\operatorname{Re}\left\{\operatorname{tr}\left(W_{1} E_{s} E_{s}^{T} X E_{s} E_{s}^{T}\right)\right\} \\
& =\operatorname{Re}\left\{\operatorname{tr}\left(E_{s} E_{s}^{T} W_{1} E_{s} E_{s}^{T} X\right)\right\}=0
\end{aligned}
$$

Similarly, we have $\left\langle E_{t} E_{t}^{T} X E_{t} E_{t}^{T}, W_{2}\right\rangle=0$ and $\left\langle E_{u} E_{u}^{T} X E_{u} E_{u}^{T}, W_{3}\right\rangle=0$. The proof is completed.

Lemma 2.8 ([15]). Let $\mathfrak{X}$ be a finite dimensional inner product space and let $\mathfrak{G}$ be the subspace of $\mathfrak{X}$. Assume that $\mathfrak{G}^{\perp}$ is the orthogonal complement of $\mathfrak{G}$. Then, for any $X \in \mathfrak{X}$, there exists an unique $G_{0} \in \mathfrak{G}$ such that

$$
\left\|X-G_{0}\right\| \leq\|X-G\|, \quad \forall G \in \mathfrak{G} .
$$

In this case, we say $G_{0}$ is the projection of $X$ onto the subspace $\mathfrak{G}$, denoted by $G_{0}=\mathcal{P}_{\mathfrak{G}}(X)$. Moreover, $G_{0}=\mathcal{P}_{\mathfrak{G}}(X)$ if and only if $\left(X-G_{0}\right) \perp \mathfrak{G}$, i.e., $\left(X-G_{0}\right) \in \mathfrak{G}^{\perp}$.
Lemma 2.9 ([27]). Let $f(X)$ be a continuous and differentiable function on a linear subspace $\mathcal{L}$ of $\mathbb{C}^{n \times n}$. Then there exists $X_{*} \in \mathcal{L}$ such that $f\left(X_{*}\right)=\min _{X \in \mathcal{L}} f(X)$ if and only if $\mathcal{P}_{\mathcal{L}}\left(\nabla f\left(X_{*}\right)\right)=O$.

Theorem 2.10. Let $R=\widetilde{Z}-\widetilde{A} X \Lambda^{2}-\widetilde{B} X \Lambda-\widetilde{C} X$. Then $\left(\widetilde{A}_{*}, \widetilde{B}_{*}, \widetilde{C}_{*}\right)$ is a solution of Problem 2.1 if and only if it is a solution of the following system of matrix equations

$$
\left\{\begin{array}{l}
R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}-S_{n}\left(R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}\right) S_{n}  \tag{2.4}\\
-E_{s} E_{s}^{T}\left[R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}-S_{n}\left(R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}\right) S_{n}\right] E_{s} E_{s}^{T}=O \\
R(X \Lambda)^{H}+X \Lambda R^{H}-S_{n}\left(R(X \Lambda)^{H}+X \Lambda R^{H}\right) S_{n} \\
-E_{t} E_{t}^{T}\left[R(X \Lambda)^{H}+X \Lambda R^{H}-S_{n}\left(R(X \Lambda)^{H}+X \Lambda R^{H}\right) S_{n}\right] E_{t} E_{t}^{T}=O \\
R X^{H}+X R^{H}-S_{n}\left(R X^{H}+X R^{H}\right) S_{n} \\
-E_{u} E_{u}^{T}\left[R X^{H}+X R^{H}-S_{n}\left(R X^{H}+X R^{H}\right) S_{n}\right] E_{u} E_{u}^{T}=O
\end{array}\right.
$$

Proof. We first define the following function on $\widetilde{S}_{1} \times \widetilde{S}_{2} \times \widetilde{S}_{3}$ :

$$
f(\widetilde{A}, \widetilde{B}, \widetilde{C})=\left\|\widetilde{Z}-\widetilde{A} X \Lambda^{2}-\widetilde{B} X \Lambda-\widetilde{C} X\right\|^{2}
$$

It is easy to see that $\left(\widetilde{A}_{*}, \widetilde{B}_{*}, \widetilde{C}_{*}\right)$ is a solution of Problem 2.1 if and only if

$$
\begin{equation*}
f\left(\widetilde{A}_{*}, \widetilde{B}_{*}, \widetilde{C}_{*}\right)=\min _{(\widetilde{A}, \widetilde{B}, \widetilde{C}) \in \widetilde{S}_{1} \times \widetilde{S}_{2} \times \widetilde{S}_{3}} f(\widetilde{A}, \widetilde{B}, \widetilde{C}) \tag{2.5}
\end{equation*}
$$

Lemma 2.9 implies that (2.5) is equivalent to

$$
\begin{equation*}
\left.\mathcal{P}_{\widetilde{S}_{1}}\left(\frac{\partial f}{\partial \widetilde{A}}\right)\right|_{\widetilde{A}=\widetilde{A}^{*}}=O,\left.\quad \mathcal{P}_{\widetilde{S}_{1}}\left(\frac{\partial f}{\partial \widetilde{B}}\right)\right|_{\widetilde{B}=\widetilde{B}^{*}}=O,\left.\quad \mathcal{P}_{\widetilde{S}_{1}}\left(\frac{\partial f}{\partial \widetilde{C}}\right)\right|_{\widetilde{C}=\widetilde{C}^{*}}=O \tag{2.6}
\end{equation*}
$$

Now we derive $\frac{\partial f}{\partial \widetilde{A}}, \frac{\partial f}{\partial \widetilde{B}}$ and $\frac{\partial f}{\partial \widetilde{C}}$. For any $\left(M_{1}, M_{1}, M_{3}\right) \in \mathbb{C}^{n \times n} \times \mathbb{C}^{n \times n} \times \mathbb{C}^{n \times n}$ and $\xi \in \mathbb{R}$, by the Taylor series expansion, we have
$f\left(\widetilde{A}+\xi M_{1}, \widetilde{B}+\xi M_{2}, \widetilde{C}+\xi M_{3}\right)=f(\widetilde{A}, \widetilde{B}, \widetilde{C})+\xi\left\langle\frac{\partial f}{\partial \widetilde{A}}, M_{1}\right\rangle+\xi\left\langle\frac{\partial f}{\partial \widetilde{B}}, M_{2}\right\rangle+\xi\left\langle\frac{\partial f}{\partial \widetilde{C}}, M_{3}\right\rangle+o(\xi)$.
According to the basic properties of Frobenius norm and matrix inner product, we get the expression

$$
\begin{aligned}
f & \left(\widetilde{A}+\xi M_{1}, \widetilde{B}+\xi M_{2}, \widetilde{C}+\xi M_{3}\right) \\
= & \left\|\widetilde{Z}-\left(\widetilde{A}+\xi M_{1}\right) X \Lambda^{2}-\left(\widetilde{B}+\xi M_{2}\right) X \Lambda-\left(\widetilde{C}+\xi M_{3}\right) X\right\|^{2} \\
= & \left\|\widetilde{Z}-\widetilde{A} X \Lambda^{2}-\widetilde{B} X \Lambda-\widetilde{C} X\right\|^{2}-2 \xi\left\langle\widetilde{Z}-\widetilde{A} X \Lambda^{2}-\widetilde{B} X \Lambda-\widetilde{C} X, M_{1} X \Lambda^{2}+M_{2} X \Lambda+M_{3} X\right\rangle \\
& +\xi^{2}\left\|M_{1} X \Lambda^{2}+M_{2} X \Lambda+M_{3} X\right\|^{2} \\
= & f(\widetilde{A}, \widetilde{B}, \widetilde{C})-2 \xi\left\langle R, M_{1} X \Lambda^{2}+M_{2} X \Lambda+M_{3} X\right\rangle+\xi^{2}\left\|M_{1} X \Lambda^{2}+M_{2} X \Lambda+M_{3} X\right\|^{2} \\
= & f(\widetilde{A}, \widetilde{B}, \widetilde{C})-2 \xi\left\langle R\left(X \Lambda^{2}\right)^{H}, M_{1}\right\rangle-2 \xi\left\langle R(X \Lambda)^{H}, M_{2}\right\rangle-2 \xi\left\langle R X^{H}, M_{3}\right\rangle+o(\xi)
\end{aligned}
$$

which, together with (2.7), yields

$$
\begin{equation*}
\frac{\partial f}{\partial \widetilde{A}}=-2 R\left(X \Lambda^{2}\right)^{H}, \frac{\partial f}{\partial \widetilde{B}}=-2 R(X \Lambda)^{H}, \frac{\partial f}{\partial \widetilde{C}}=-2 R X^{H} \tag{2.8}
\end{equation*}
$$

For any $T_{1} \in \widetilde{S}_{1}$, by Lemmas 2.6, 2.7 and 2.8 , we obtain

$$
\begin{align*}
& \left\langle-2 R\left(X \Lambda^{2}\right)^{H}, T_{1}\right\rangle=-\frac{1}{2}\left\langle R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}-S_{n}\left(R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}\right) S_{n}, T_{1}\right\rangle \\
= & -\frac{1}{2}\left\langle R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}-S_{n}\left(R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}\right) S_{n}-E_{s} E_{s}^{T}\left[R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}\right.\right. \\
& \left.\left.-S_{n}\left(R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}\right) S_{n}\right] E_{s} E_{s}^{T}, T_{1}\right\rangle . \tag{2.9}
\end{align*}
$$

and

$$
\begin{align*}
\left\langle-2 R\left(X \Lambda^{2}\right)^{H}, T_{1}\right\rangle & =\left\langle-2 R\left(X \Lambda^{2}\right)^{H}-\mathcal{P}_{\widetilde{S}_{1}}\left(-2 R\left(X \Lambda^{2}\right)^{H}\right), T_{1}\right\rangle+\left\langle\mathcal{P}_{\widetilde{S}_{1}}\left(-2 R\left(X \Lambda^{2}\right)^{H}\right), T_{1}\right\rangle \\
& =\left\langle\mathcal{P}_{\widetilde{S}_{1}}\left(-2 R\left(X \Lambda^{2}\right)^{H}\right), T_{1}\right\rangle \tag{2.10}
\end{align*}
$$

Together (2.9) and (2.10) yields

$$
\begin{aligned}
& \left\langle\mathcal{P}_{\widetilde{S}_{1}}\left(-2 R\left(X \Lambda^{2}\right)^{H}\right), T_{1}\right\rangle=-\frac{1}{2}\left\langle R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}-S_{n}\left(R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}\right) S_{n}\right. \\
& \left.-E_{s} E_{s}^{T}\left[R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}-S_{n}\left(R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}\right) S_{n}\right] E_{s} E_{s}^{T}, T_{1}\right\rangle
\end{aligned}
$$

It follows from the relation (2.8) and the arbitrariness of $T_{1}$ that

$$
\begin{aligned}
\mathcal{P}_{\widetilde{S}_{1}}\left(\frac{\partial f}{\partial \widetilde{A}}\right)= & \mathcal{P}_{\widetilde{S}_{1}}\left(-2 R\left(X \Lambda^{2}\right)^{H}\right) \\
= & -\frac{1}{2}\left\{R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}+S_{n}\left(R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}\right) S_{n}\right. \\
& \left.\quad-E_{s} E_{s}^{T}\left[R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}+S_{n}\left(R\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}\right) S_{n}\right] E_{s} E_{s}^{T}\right\} .
\end{aligned}
$$

Similarly

$$
\left\{\begin{array}{l}
\mathcal{P}_{\widetilde{S}_{2}}\left(\frac{\partial f}{\partial \widetilde{B}}\right)=-\mathcal{P}_{\widetilde{S}_{2}}\left(2 R(X \Lambda)^{H}\right)=-\frac{1}{2}\left\{R(X \Lambda)^{H}+X \Lambda R^{H}+S_{n}\left(R(X \Lambda)^{H}+X \Lambda R^{H}\right) S_{n}\right. \\
\left.\quad-E_{t} E_{t}^{T}\left[R(X \Lambda)^{H}+X \Lambda R^{H}+S_{n}\left(R(X \Lambda)^{H}+X \Lambda R^{H}\right) S_{n}\right] E_{t} E_{t}^{T}\right\} \\
\mathcal{P}_{\widetilde{S}_{3}}\left(\frac{\partial f}{\partial \widetilde{C}}\right)=-\mathcal{P}_{\widetilde{S}_{3}}\left(2 R X^{H}\right)=-\frac{1}{2}\left\{R X^{H}+X R^{H}+S_{n}\left(R X^{H}+X R^{H}\right) S_{n}\right. \\
\left.\quad-E_{u} E_{u}^{T}\left[R X^{H}+X R^{H}+S_{n}\left(R X^{H}+X R^{H}\right) S_{n}\right] E_{u} E_{u}^{T}\right\}
\end{array}\right.
$$

By the relation (2.6), the result is established. The proof is completed.
It is well known that many iterative methods for matrix equations are based on iterative methods for linear system $A x=b$, where $A \in \mathbb{C}^{m \times n}$ and $b \in \mathbb{C}^{m}$. For the least squares problem

$$
\begin{equation*}
\min _{x \in \mathbb{C}^{n}}\|b-A x\| \text { for } A \in \mathbb{C}^{m \times n} \text { and } b \in \mathbb{C}^{m} \tag{2.11}
\end{equation*}
$$

the LSQR and CGLS methods [8, 34] are all efficient. More importantly, CG-type method overcomes the shortcoming of slow convergence of steepest descent method, and avoids storing and computing Hessian matrix and its inverse of Newton method. Many researchers have investigations into this direction. For instance, Zhou et al. [42] established a gradientbased iterative algorithm for solving the coupled matrix equation. Based on CG method, Lv and Zhang [32] raised a neat iterative algorithm in the least squares sense for solving a periodic Sylvester matrix equation.

Motivated by the work mentioned above, we propose a nonlinear conjugate gradient method for solving Problem 2.1.

Algorithm 2.11. (Nonlinear conjugate gradient method (NCG) for Problem 2.1)
Step 0 Input $X \in \mathbb{C}^{n \times m}, \Lambda=\operatorname{diag}\left(\lambda_{1}, \lambda_{2}, \ldots, \lambda_{m}\right) \in \mathbb{C}^{m \times m}, s=\left(s_{1}, s_{2}, \ldots, s_{p}, n+1-\right.$ $\left.s_{p}, \ldots, n+1-s_{2}, n+1-s_{1}\right) \in D_{2 p, n}, t=\left(t_{1}, t_{2}, \ldots, t_{q}, n+1-t_{q}, \ldots, n+1-t_{2}, n+1-t_{1}\right) \in$ $D_{2 q, n}, u=\left(u_{1}, u_{2}, \ldots, u_{r}, n+1-u_{r}, \ldots, n+1-u_{2}, n+1-u_{1}\right) \in D_{2 r, n}, A_{p} \in \mathbb{C}^{2 p \times 2 p}, B_{q} \in$ $\mathbb{C}^{2 q \times 2 q}, C_{r} \in \mathbb{C}^{2 r \times 2 r}, \widetilde{A}_{p} \in \mathbb{C}^{n \times n}, \widetilde{B}_{q} \in \mathbb{C}^{n \times n}, \widetilde{C}_{r} \in \mathbb{C}^{n \times n}$ and $\widetilde{Z}=-\widetilde{A}_{p} X \Lambda^{2}-\widetilde{B}_{q} X \Lambda-\widetilde{C}_{r} X$. Choose the initial matrices $\widetilde{A}(1) \in \widetilde{S}_{1}, \widetilde{B}(1) \in \widetilde{S}_{2}$ and $\widetilde{C}(1) \in \widetilde{S}_{3}$.

Step 1 Compute

$$
\begin{aligned}
R(1)= & \widetilde{Z}-\widetilde{A}(1) X \Lambda^{2}-\widetilde{B}(1) X \Lambda-\widetilde{C}(1) X, \\
P_{1}(1)= & R(1)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(1)-S_{n}\left(R(1)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(1)\right) S_{n} \\
& -E_{s} E_{s}^{T}\left[R(1)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(1)-S_{n}\left(R(1)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(1)\right) S_{n}\right] E_{s} E_{s}^{T}, \\
P_{2}(1)= & R(1)(X \Lambda)^{H}+X \Lambda R^{H}(1)-S_{n}\left(R(1)(X \Lambda)^{H}+X \Lambda R^{H}(1)\right) S_{n} \\
& -E_{t} E_{t}^{T}\left[R(1)(X \Lambda)^{H}+X \Lambda R^{H}(1)-S_{n}\left(R(1)(X \Lambda)^{H}+X \Lambda R^{H}(1)\right) S_{n}\right] E_{t} E_{t}^{T}, \\
P_{3}(1)= & R(1) X^{H}+X R(1)^{H}-S_{n}\left(R(1) X^{H}+X R^{H}(1)\right) S_{n} \\
& -E_{u} E_{u}^{T}\left[R(1) X^{H}+X R^{H}(1)-S_{n}\left(R(1) X^{H}+X R^{H}(1)\right) S_{n}\right] E_{u} E_{u}^{T} .
\end{aligned}
$$

Set $M_{1}(1)=-P_{1}(1), M_{2}(1)=-P_{2}(1)$ and $M_{3}(1)=-P_{3}(1)$. Denote

$$
P(1)=\left(\begin{array}{ccc}
P_{1}(1) & O & O \\
O & P_{2}(1) & O \\
O & O & P_{3}(1)
\end{array}\right) \quad \text { and } \quad M(1)=\left(\begin{array}{ccc}
M_{1}(1) & O & O \\
O & M_{2}(1) & O \\
O & O & M_{3}(1)
\end{array}\right)
$$

Set $k=1$.

Step 2 If $\|R(k)\|=0$ or $\|P(k)\|=0$, stop. Otherwise, go to Step 3.
Step 3 Compute $Q(k)=M_{1}(k) X \Lambda^{2}+M_{2}(k) X \Lambda+M_{3}(k) X$ and

$$
\begin{equation*}
\alpha_{k}=\frac{\langle M(k), P(k)\rangle}{4\|Q(k)\|^{2}} \tag{2.12}
\end{equation*}
$$

Update the sequences

$$
\begin{aligned}
& \widetilde{A}(k+1)=\widetilde{A}(k)+\alpha_{k} M_{1}(k), \\
& \widetilde{B}(k+1)=\widetilde{B}(k)+\alpha_{k} M_{2}(k), \\
& \widetilde{C}(k+1)=\widetilde{C}(k)+\alpha_{k} M_{3}(k), R(k+1)=R(k)-\alpha_{k} Q(k), \\
& P_{1}(k+1)= R(k+1)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(k+1) \\
&-S_{n}\left(R(k+1)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(k+1)\right) S_{n} \\
&-E_{s} E_{s}^{T}\left[R(k+1)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(k+1)\right. \\
&\left.-S_{n}\left(R(k+1)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(k+1)\right) S_{n}\right] E_{s} E_{s}^{T} \\
& P_{2}(k+1)= R(k+1)(X \Lambda)^{H}+X \Lambda R^{H}(k+1)-S_{n}\left(R(k+1)(X \Lambda)^{H}\right. \\
&\left.+X \Lambda R^{H}(k+1)\right) S_{n}-E_{t} E_{t}^{T}\left[R(k+1)(X \Lambda)^{H}+X \Lambda R^{H}(k+1)\right. \\
&\left.-S_{n}\left(R(k+1)(X \Lambda)^{H}+X \Lambda R^{H}(k+1)\right) S_{n}\right] E_{t} E_{t}^{T} \\
& P_{3}(k+1)= R(k+1) X^{H}+X R(k+1)^{H}-S_{n}\left(R(k+1) X^{H}\right. \\
&\left.+X R^{H}(k+1)\right) S_{n}-E_{u} E_{u}^{T}\left[R(k+1) X^{H}+X R^{H}(k+1)\right. \\
&\left.-S_{n}\left(R(k+1) X^{H}+X R^{H}(k+1)\right) S_{n}\right] E_{u} E_{u}^{T}
\end{aligned}
$$

Denote

$$
P(k+1)=\left(\begin{array}{ccc}
P_{1}(k+1) & O & O \\
O & P_{2}(k+1) & O \\
O & O & P_{3}(k+1)
\end{array}\right)
$$

Compute

$$
\begin{equation*}
\beta_{k}=\frac{\|P(k+1)\|^{2}}{\|P(k)\|^{2}} \tag{2.13}
\end{equation*}
$$

Update the sequences

$$
M_{j}(k+1)=-P_{j}(k+1)+\beta_{k} M_{j}(k)
$$

for $j=1,2,3$. Denote

$$
M(k+1)=\left(\begin{array}{ccc}
M_{1}(k+1) & O & O \\
O & M_{2}(k+1) & O \\
O & O & M_{3}(k+1)
\end{array}\right)
$$

Step 4 Set $k=k+1$, go to Step 2 .
Remark 2.12. According to Lemma 2.1 and Algorithm 2.11, it follows that matrix sequences generated by Algorithm 2.11 satisfy $\{\widetilde{A}(k)\},\left\{P_{1}(k)\right\},\left\{M_{1}(k)\right\} \subseteq \widetilde{S}_{1},\{\widetilde{B}(k)\}$, $\left\{P_{2}(k)\right\},\left\{M_{2}(k)\right\} \subseteq \widetilde{S}_{2}$ and $\{\widetilde{C}(k)\},\left\{P_{3}(k)\right\},\left\{M_{3}(k)\right\} \subseteq \widetilde{S}_{3}$.

Remark 2.13. If $\|R(k)\|=0$, then $(\widetilde{A}(k), \widetilde{B}(k), \widetilde{C}(k))$ is a solution of matrix equation $\widetilde{A} X \Lambda^{2}+\widetilde{B} X \Lambda+\widetilde{C} X-\widetilde{Z}=O$. That is, $(\widetilde{A}(k), \widetilde{B}(k), \widetilde{C}(k))$ is a solution of Problem 2.1 and

$$
\min _{(\widetilde{A}, \widetilde{B}, \widetilde{C}) \in \widetilde{S}_{1} \times \widetilde{S}_{2} \times \widetilde{S}_{3}}\left\|\widetilde{A} X \Lambda^{2}+\widetilde{B} X \Lambda+\widetilde{C} X-\widetilde{Z}\right\|=0
$$

If $\|P(k)\|=0$, then $(\widetilde{A}(k), \widetilde{B}(k), \widetilde{C}(k))$ is a solution of Problem 2.1 and

$$
\min _{(\widetilde{A}, \widetilde{B}, \widetilde{C}) \in \widetilde{S}_{1} \times \widetilde{S}_{2} \times \widetilde{S}_{3}}\left\|\widetilde{A} X \Lambda^{2}+\widetilde{B} X \Lambda+\widetilde{C} X-\widetilde{Z}\right\|>0
$$

Now, we list some basic properties of Algorithm 2.11.
Lemma 2.14. Let the sequences $\{P(k)\}$ and $\{M(k)\}$ be generated by Algorithm 2.11. Then

$$
\langle M(k), P(k+1)\rangle=0 .
$$

Proof. We first define $G_{1}(k), G_{2}(k)$ and $G_{3}(k)$ as follows:

$$
\left\{\begin{array}{c}
G_{1}(k)=Q(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Q^{H}(k)-S_{n}\left(Q(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Q^{H}(k)\right) S_{n} \\
\quad-E_{s} E_{s}^{T}\left[Q(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(k)-S_{n}\left(Q(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Q^{H}(k)\right) S_{n}\right] E_{s} E_{s}^{T} \\
G_{2}(k)=Q(k)(X \Lambda)^{H}+X \Lambda Q^{H}(k)-S_{n}\left(Q(k)(X \Lambda)^{H}+X \Lambda Q^{H}(k)\right) S_{n} \\
\quad-E_{t} E_{t}^{T}\left[Q(k)(X \Lambda)^{H}+X \Lambda R^{H}(1)-S_{n}\left(Q(1)(X \Lambda)^{H}+X \Lambda Q^{H}(1)\right) S_{n}\right] E_{t} E_{t}^{T} \\
G_{3}(k)=Q(1) X^{H}+X Q(k)^{H}-S_{n}\left(Q(k) X^{H}+X Q^{H}(k)\right) S_{n} \\
-E_{u} E_{u}^{T}\left[Q(k) X^{H}+X Q^{H}(k)-S_{n}\left(Q(k) X^{H}+X Q^{H}(k)\right) S_{n}\right] E_{u} E_{u}^{T}
\end{array}\right.
$$

Since $R(k+1)=R(k)-\alpha_{k} Q(k)$, it follows that

$$
\begin{aligned}
P_{1}(k+1)= & R(k+1)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(k+1)-S_{n}\left(R(k+1)\left(X \Lambda^{2}\right)^{H}\right. \\
& \left.+X \Lambda^{2} R^{H}(k+1)\right) S_{n}-E_{s} E_{s}^{T}\left[R(k+1)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(k+1)\right. \\
& \left.-S_{n}\left(R(k+1)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(k+1)\right) S_{n}\right] E_{s} E_{s}^{T} \\
= & P_{1}(k)-\alpha_{k}\left[Q(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Q^{H}(k)-S_{n}\left(Q(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Q^{H}(k)\right) S_{n}\right. \\
& -E_{s} E_{s}^{T}\left[Q(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(k)-S_{n}\left(Q(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Q^{H}(k)\right) S_{n}\right] E_{s} E_{s}^{T} \\
= & P_{1}(k)-\alpha_{k} G_{1}(k) .
\end{aligned}
$$

Similarly, we have $P_{2}(k+1)=P_{2}(k)-\alpha_{k} G_{2}(k)$ and $P_{3}(k+1)=P_{3}(k)-\alpha_{k} G_{3}(k)$. Then

$$
\langle M(k), P(k+1)\rangle=\sum_{j=1}^{3}\left\langle M_{j}(k), P_{j}(k)-\alpha_{k} G_{j}(k)\right\rangle=\langle M(k), P(k)\rangle-\alpha_{k}\langle M(k), G(k)\rangle .(
$$

On the other hand, by Algorithm 2.11, Lemmas 2.1, 2.6 and 2.7 and Remark 2.12, it follows that

$$
\begin{aligned}
\left\langle M_{1}(k), G_{1}(k)\right\rangle= & \left\langle M_{1}(k), Q(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Q^{H}(k)\right. \\
& -S_{n}\left(Q(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Q^{H}(k)\right) S_{n} \\
& -E_{s} E_{s}^{T}\left[Q(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Q^{H}(k)\right. \\
& \left.\left.-S_{n}\left(Q(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Q^{H}(k)\right) S_{n}\right] E_{s} E_{s}^{T}\right\rangle \\
= & \left\langle M_{1}(k), Q(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Q^{H}(k)\right. \\
& \left.-S_{n}\left(Q(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Q^{H}(k)\right) S_{n}\right\rangle \\
= & 4\left\langle M_{1}(k), Q(k)\left(X \Lambda^{2}\right)^{H}\right\rangle=4\left\langle M_{1}(k) X \Lambda^{2}, Q(k)\right\rangle .
\end{aligned}
$$

Similarly, we have

$$
\left\langle M_{2}(k), G_{2}(k)\right\rangle=4\left\langle M_{2}(k) X \Lambda, Q(k)\right\rangle \text { and }\left\langle M_{3}(k), G_{3}(k)\right\rangle=4\left\langle M_{3}(k) X, Q(k)\right\rangle
$$

Hence

$$
\langle M(k), G(k)\rangle=\sum_{j=1}^{3}\left\langle M_{j}(k), G_{j}(k)\right\rangle=4\langle Q(k), Q(k)\rangle
$$

which, together with (2.12) and (2.14), yields

$$
\begin{aligned}
\langle M(k), P(k+1)\rangle & =\langle M(k), P(k)\rangle-\alpha_{k}\langle M(k), G(k)\rangle=\langle M(k), P(k)\rangle-4 \alpha_{k}\langle Q(k), Q(k)\rangle \\
& =\langle M(k), P(k)\rangle-4 \times \frac{\langle M(k), P(k)\rangle}{4\|Q(k)\|^{2}}\|Q(k)\|^{2}=0 .
\end{aligned}
$$

The proof is completed.
Lemma 2.15. Let the sequences $\{M(k)\}$ and $\{P(k)\}$ be generated by Algorithm 2.11. Then

$$
\langle M(k), P(k)\rangle=-\|P(k)\|^{2}
$$

Proof. For $k=1$, since $M_{j}(1)=-P_{j}(1)$, we immediately have

$$
\langle M(1), P(1)\rangle=\sum_{j=1}^{3}\left\langle M_{j}(1), P_{j}(1)\right\rangle=-\sum_{j=1}^{3}\left\langle P_{j}(1), P_{j}(1)\right\rangle=-\|P(1)\|^{2}
$$

For $k \geq 2$, it follows from Lemma 2.14 that

$$
\begin{aligned}
\langle M(k+1), P(k+1)\rangle & =\sum_{j=1}^{3}\left\langle-P_{j}(k+1)+\beta_{k} M_{j}(k), P_{j}(k+1)\right\rangle \\
& =-\|P(k+1)\|^{2}+\beta_{k}\langle M(k), P(k+1)\rangle=-\|P(k+1)\|^{2}
\end{aligned}
$$

The proof is completed.
Lemma 2.16. Let the sequences $\{M(k)\}$ and $\{P(k)\}$ be generated by Algorithm 2.11. Then

$$
\begin{equation*}
\sum_{k=1}^{\infty} \frac{\|P(k)\|^{4}}{\|M(k)\|^{2}}<\infty \tag{2.15}
\end{equation*}
$$

Proof. By direct calculations, we have

$$
\begin{align*}
& \|Q(k)\|^{2}=\left\|M_{1}(k) X \Lambda^{2}+M_{2}(k) X \Lambda+M_{3}(k) X\right\|^{2} \\
& =\left\|\left(\left(X \Lambda^{2}\right)^{T} \otimes I\right) \operatorname{vec}\left(M_{1}(k)\right)+\left((X \Lambda)^{T} \otimes I\right) \operatorname{vec}\left(M_{2}(k)\right)+\left(X^{T} \otimes I\right) \operatorname{vec}\left(M_{3}(k)\right)\right\|^{2} \\
& =\left\|\left(\begin{array}{lll}
\left(X \Lambda^{2}\right)^{T} \otimes I & (X \Lambda)^{T} \otimes I & X^{T} \otimes I
\end{array}\right)\left(\begin{array}{c}
\operatorname{vec}\left(M_{1}(k)\right) \\
\operatorname{vec}\left(M_{2}(k)\right) \\
\operatorname{vec}\left(M_{3}(k)\right)
\end{array}\right)\right\|^{2} \\
& \leq \|\left(\begin{array}{lll}
\left(X \Lambda^{2}\right)^{T} \otimes I & (X \Lambda)^{T} \otimes I & \left.X^{T} \otimes I\right)\left\|^{2}\right\|\left(\begin{array}{c}
\operatorname{vec}\left(M_{1}(k)\right) \\
\operatorname{vec}\left(M_{2}(k)\right) \\
\operatorname{vec}\left(M_{3}(k)\right)
\end{array}\right) \|^{2}
\end{array} \|^{2}\right. \\
& =\theta \sum_{j=1}^{3}\left\|\operatorname{vec}\left(M_{j}(k)\right)\right\|^{2}=\theta \sum_{j=1}^{3}\left\|M_{j}(k)\right\|^{2}=\theta\|M(k)\|^{2}, \tag{2.16}
\end{align*}
$$

where

$$
\theta=\|\left(\begin{array}{ll}
\left.\left(X \Lambda^{2}\right)^{T} \otimes I \quad(X \Lambda)^{T} \otimes I \quad X^{T} \otimes I\right)
\end{array} \|^{2} .\right.
$$

According to Algorithm 2.11, Lemmas 2.6 and 2.7, we obtain

$$
\begin{aligned}
\| & R(k+1)\left\|^{2}=\right\| R(k)-\alpha_{k} Q(k)\left\|^{2}=\right\| R(k)\left\|^{2}-2 \alpha_{k}\langle R(k), Q(k)\rangle+\alpha_{k}^{2}\right\| Q(k) \|^{2} \\
= & \|R(k)\|^{2}-2 \alpha_{k}\left\langle R(k), M_{1}(k) X \Lambda^{2}+M_{2}(k) X \Lambda+M_{3}(k) X\right\rangle+\alpha_{k}^{2}\|Q(k)\|^{2} \\
= & \|R(k)\|^{2}-2 \alpha_{k}\left[\left\langle R(k)\left(X \Lambda^{2}\right)^{H}, M_{1}(k)\right\rangle+\left\langle R(k)(X \Lambda)^{H}, M_{2}(k)\right\rangle+\left\langle R(k) X^{H}, M_{3}(k)\right\rangle\right] \\
& +\alpha_{k}^{2}\|Q(k)\|^{2} \\
= & \|R(k)\|^{2}-\frac{\alpha_{k}}{2}\left\{\left\langleR(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(k)-S_{n}\left(R(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(k)\right) S_{n}\right.\right. \\
& \left.-E_{s} E_{s}^{T}\left[R(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(k)-S_{n}\left(R(k)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R^{H}(k)\right) S_{n}\right] E_{s} E_{s}^{T}, M_{1}(k)\right\rangle \\
& +\left\langle R(k)(X \Lambda)^{H}+X \Lambda R^{H}(k)-S_{n}\left(R(k)(X \Lambda)^{H}+X \Lambda R^{H}(k)\right) S_{n}\right. \\
& \left.-E_{t} E_{t}^{T}\left[R(k)(X \Lambda)^{H}+X \Lambda R^{H}(k)-S_{n}\left(R(k)(X \Lambda)^{H}+X \Lambda R^{H}(k)\right) S_{n}\right] E_{t} E_{t}^{T}, M_{2}(k)\right\rangle \\
& +\left\langle R(k) X^{H}+X R^{H}(k)-S_{n}\left(R(k) X^{H}+X R^{H}(k)\right) S_{n}\right. \\
& \left.\left.-E_{u} E_{u}^{T}\left[R(k) X^{H}+X R^{H}(k)-S_{n}\left(R(k) X^{H}+X R^{H}(k)\right) S_{n}\right] E_{u} E_{u}^{T}, M_{3}(k)\right\rangle\right\} \\
& +\alpha_{k}^{2}\|Q(k)\|^{2} \\
= & \|R(k)\|^{2}-\frac{\alpha_{k}}{2}\left[\left\langle P_{1}(k), M_{1}(k)\right\rangle+\left\langle P_{2}(k), M_{2}(k)\right\rangle+\left\langle P_{3}(k), M_{3}(k)\right\rangle\right]+\alpha_{k}^{2}\|Q(k)\|^{2} \\
= & \|R(k)\|^{2}-\frac{\alpha_{k}}{2}\langle P(k), M(k)\rangle+\alpha_{k} \frac{\langle M(k), P(k)\rangle}{4\|Q(k)\|^{2}}\|Q(k)\|^{2} \\
= & \|R(k)\|^{2}-\frac{\alpha_{k}}{4}\langle P(k), M(k)\rangle .
\end{aligned}
$$

This implies that

$$
\begin{equation*}
\|R(k+1)\|^{2}-\|R(k)\|^{2}=-\frac{\alpha_{k}}{4}\langle P(k), M(k)\rangle=-\frac{1}{16} \frac{\langle M(k), P(k)\rangle^{2}}{\|Q(k)\|^{2}} \leq 0 . \tag{2.17}
\end{equation*}
$$

Then the sequence $\left\{\|R(k)\|^{2}\right\}$ is decreasing and its limit exists. Hence, by Lemma 2.15 and the relations (2.16) and (2.17), it follows that

$$
\begin{aligned}
\sum_{k=1}^{\infty} \frac{\|P(k)\|^{4}}{\|M(k)\|^{2}} & =\sum_{k=1}^{\infty} \frac{\langle M(k), P(k)\rangle^{2}}{\|M(k)\|^{2}} \leq \sum_{k=1}^{\infty} \theta \frac{\langle M(k), P(k)\rangle^{2}}{\|Q(k)\|^{2}} \\
& =16 \theta \sum_{k=1}^{\infty}\left[\|R(k)\|^{2}-\|R(k+1)\|^{2}\right]=16 \theta\left[\|R(1)\|^{2}-\lim _{k \rightarrow \infty}\|R(k)\|^{2}\right]<\infty
\end{aligned}
$$

The proof is completed.

Theorem 2.17. Let $\left\{P_{j}(k)\right\}(j=1,2,3)$ be generated by Algorithm 2.11. Then $\lim _{k \rightarrow \infty}\left\|P_{j}(k)\right\|=0$ for $j=1,2,3$.

Proof. By Lemma 2.14, it follows that

$$
\begin{aligned}
& \|M(k+1)\|^{2}=\sum_{j=1}^{3}\left\|M_{j}(k+1)\right\|^{2}=\sum_{j=1}^{3}\left\|-P_{j}(k+1)+\beta_{k} M_{j}(k)\right\|^{2} \\
= & \|P(k+1)\|^{2}-2 \beta_{k}\langle P(k+1), M(k)\rangle+\beta_{k}^{2}\|M(k)\|^{2} \\
= & \|P(k+1)\|^{2}+\beta_{k}^{2}\|M(k)\|^{2}=\|P(k+1)\|^{2}+\frac{\|P(k+1)\|^{4}}{\|P(k)\|^{4}}\|M(k)\|^{2} .
\end{aligned}
$$

Then

$$
\begin{equation*}
\frac{\|M(k+1)\|^{2}}{\|P(k+1)\|^{4}}=\frac{1}{\|P(k+1)\|^{2}}+\frac{\|M(k)\|^{2}}{\|P(k)\|^{4}} \tag{2.18}
\end{equation*}
$$

Let $t_{k}=\frac{\|M(k)\|^{2}}{\|P(k)\|^{4}}$. According to the relation (2.18), we immediately have

$$
\begin{equation*}
t_{k+1}=t_{k}+\frac{1}{\|P(k+1)\|^{2}} \tag{2.19}
\end{equation*}
$$

We now prove this result by contradiction. Assume that $\lim _{k \rightarrow \infty}\|P(k)\| \neq 0$. Then there exists a positive number $\delta>0$ such that $\|P(k)\| \geq \delta$ for all $k \geq 1$. It follows from the relation (2.19) that

$$
t_{k+1} \leq t_{0}+\frac{k+1}{\delta}
$$

which implies that

$$
\begin{equation*}
\sum_{k=1}^{\infty} \frac{1}{t_{k}} \geq \sum_{k=1}^{\infty} \frac{\delta}{\delta t_{0}+k+1}=\infty \tag{2.20}
\end{equation*}
$$

However, by Lemma 2.16, we get

$$
\sum_{k=1}^{\infty} \frac{1}{t_{k}}=\sum_{k=1}^{\infty} \frac{\|P(k)\|^{4}}{\|M(k)\|^{2}}<\infty
$$

which obtains a contradiction. Hence $\lim _{k \rightarrow \infty}\|P(k)\|=0$, and then $\lim _{k \rightarrow \infty}\left\|P_{j}(k)\right\|=0$ for $j=1,2,3$. The proof is completed.

Finally, we consider the least Frobenius norm solution of Problem 2.1. The following lemma is needed for the main result.

Lemma 2.18. Suppose that $\left(\widetilde{A}_{*}, \widetilde{B}_{*}, \widetilde{C}_{*}\right)$ is a solution of Problem 2.1. Then, an arbitrary solution $(\widetilde{A}, \widetilde{B}, \widetilde{C})$ of Problem 2.1 can be expressed as

$$
\begin{equation*}
(\widetilde{A}, \widetilde{B}, \widetilde{C})=\left(\widetilde{A}_{*}+\mathcal{W}_{1}, \widetilde{B}_{*}+\mathcal{W}_{2}, \widetilde{C}_{*}+\mathcal{W}_{3}\right) \tag{2.21}
\end{equation*}
$$

where $\left(\mathcal{W}_{1}, \mathcal{W}_{2}, \mathcal{W}_{3}\right) \in \widetilde{S}_{1} \times \widetilde{S}_{2} \times \widetilde{S}_{3}$ satisfies

$$
\begin{equation*}
\mathcal{W}_{1} X \Lambda^{2}+\mathcal{W}_{2} X \Lambda+\mathcal{W}_{3} X=O \tag{2.22}
\end{equation*}
$$

Proof. For an arbitrary solution $(\widetilde{A}, \widetilde{B}, \widetilde{C})$ of Problem 2.1, we first define the following matrices

$$
\mathcal{W}_{1}=\widetilde{A}-\widetilde{A}_{*}, \quad \mathcal{W}_{2}=\widetilde{B}-\widetilde{B}_{*}, \quad \mathcal{W}_{3}=\widetilde{C}-\widetilde{C}_{*}
$$

Obviously $\left(\mathcal{W}_{1}, \mathcal{W}_{2}, \mathcal{W}_{3}\right) \in \widetilde{S}_{1} \times \widetilde{S}_{2} \times \widetilde{S}_{3}$. Let $R_{*}=\widetilde{A}_{*} X \Lambda^{2}+\widetilde{B}_{*} X \Lambda+\widetilde{C}_{*} X-\widetilde{Z}$. Since $\left(\widetilde{A}_{*}, \widetilde{B}_{*}, \widetilde{C}_{*}\right)$ is a solution of Problem 2.1, it follows from Theorem 2.10 that

$$
\begin{aligned}
& \left\langle\widetilde{A}_{*} X \Lambda^{2}+\widetilde{B}_{*} X \Lambda+\widetilde{C}_{*} X-\widetilde{Z}, \mathcal{W}_{1} X \Lambda^{2}+\mathcal{W}_{2} X \Lambda+\mathcal{W}_{3} X\right\rangle \\
= & \left\langle R_{*}\left(X \Lambda^{2}\right)^{H}, \mathcal{W}_{1}\right\rangle+\left\langle R_{*}(X \Lambda)^{H}, \mathcal{W}_{2}\right\rangle+\left\langle R_{*} X^{H}, \mathcal{W}_{3}\right\rangle \\
= & \frac{1}{4}\left\{\left\langleR_{*}\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R_{*}^{H}-S_{n}\left(R_{*}\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R_{*}^{H}\right) S_{n}\right.\right. \\
& \left.-E_{s} E_{s}^{T}\left[R_{*}\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R_{*}^{H}-S_{n}\left(R_{*}\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} R_{*}^{H}\right) S_{n}\right] E_{s} E_{s}^{T}, \mathcal{W}_{1}\right\rangle \\
& +\left\langle R_{*}(X \Lambda)^{H}+X \Lambda R_{*}^{H}-S_{n}\left(R_{*}(X \Lambda)^{H}+X \Lambda R_{*}^{H}\right) S_{n}\right. \\
& \left.-E_{t} E_{t}^{T}\left[R_{*}(X \Lambda)^{H}+X \Lambda R_{*}^{H}-S_{n}\left(R_{*}(X \Lambda)^{H}+X \Lambda R_{*}^{H}\right) S_{n}\right] E_{t} E_{t}^{T}, \mathcal{W}_{2}\right\rangle \\
& +\left\langle R_{*} X^{H}+X R_{*}^{H}-S_{n}\left(R_{*} X^{H}+X R_{*}^{H}\right) S_{n}\right. \\
& \left.\left.-E_{u} E_{u}^{T}\left[R_{*} X^{H}+X R_{*}^{H}-S_{n}\left(R_{*} X^{H}+X R_{*}^{H}\right) S_{n}\right] E_{u} E_{u}^{T}, \mathcal{W}_{3}\right\rangle\right\}=0 .
\end{aligned}
$$

Then

$$
\begin{aligned}
& \left\|\widetilde{A} X \Lambda^{2}+\widetilde{B} X \Lambda+\widetilde{C} X-\widetilde{Z}\right\|^{2} \\
= & \left\|\left(\widetilde{A}_{*}+\mathcal{W}_{1}\right) X \Lambda^{2}+\left(\widetilde{B}_{*}+\mathcal{W}_{2}\right) X \Lambda+\left(\widetilde{C}_{*}+\mathcal{W}_{3}\right) X-\widetilde{Z}\right\|^{2} \\
= & \left\|\widetilde{A}_{*} X \Lambda^{2}+\widetilde{B}_{*} X \Lambda+\widetilde{C}_{*} X-\widetilde{Z}\right\|^{2} \\
& +2\left\langle\widetilde{A}_{*} X \Lambda^{2}+\widetilde{B}_{*} X \Lambda+\widetilde{C}_{*} X-\widetilde{Z}, \mathcal{W}_{1} X \Lambda^{2}+\mathcal{W}_{2} X \Lambda+\mathcal{W}_{3} X\right\rangle \\
& +\left\|\mathcal{W}_{1} X \Lambda^{2}+\mathcal{W}_{2} X \Lambda+\mathcal{W}_{3} X\right\|^{2} \\
= & \left\|\widetilde{A}_{*} X \Lambda^{2}+\widetilde{B}_{*} X \Lambda+\widetilde{C}_{*} X-\widetilde{Z}\right\|^{2}+\left\|\mathcal{W}_{1} X \Lambda^{2}+\mathcal{W}_{2} X \Lambda+\mathcal{W}_{3} X\right\|^{2} .
\end{aligned}
$$

This implies that

$$
\mathcal{W}_{1} X \Lambda^{2}+\mathcal{W}_{2} X \Lambda+\mathcal{W}_{3} X=O
$$

The proof is completed.

Theorem 2.19. If we choose the initial matrices $\widetilde{A}(1), \widetilde{B}(1)$ and $\widetilde{C}(1)$ as follows:

$$
\left\{\begin{align*}
\widetilde{A}(1) & =Y(1)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Y^{H}(1)-S_{n}\left(Y(1)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Y^{H}(1)\right) S_{n}  \tag{2.23}\\
& -E_{s} E_{s}^{T}\left[Y(1)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Y^{H}(1)-S_{n}\left(Y(1)\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Y^{H}(1)\right) S_{n}\right] E_{s} E_{s}^{T}, \\
\widetilde{B}(1) & =Y(1)(X \Lambda)^{H}+X \Lambda Y^{H}(1)-S_{n}\left(Y(1)(X \Lambda)^{H}+X \Lambda Y(1)^{H}\right) S_{n} \\
& -E_{t} E_{t}^{T}\left[Y(1)(X \Lambda)^{H}+X \Lambda Y^{H}(1)-S_{n}\left(Y(1)(X \Lambda)^{H}+X \Lambda Y^{H}(1)\right) S_{n}\right] E_{t} E_{t}^{T}, \\
\widetilde{C}(1) & =Y(1) X^{H}+X Y(1)^{H}-S_{n}\left(Y(1) X^{H}+X Y^{H}(1)\right) S_{n} \\
& -E_{u} E_{u}^{T}\left[Y(1) X^{H}+X Y^{H}(1)-S_{n}\left(Y(1) X^{H}+X Y^{H}(1)\right) S_{n}\right] E_{u} E_{u}^{T},
\end{align*}\right.
$$

where $Y(1) \in \mathbb{C}^{n \times m}$ is an arbitrary matrix (especially, take $Y(1)=O$ ), then the solution $\left(\widetilde{A}_{*}, \widetilde{B}_{*}, \widetilde{C}_{*}\right)$ given by Algorithm 2.11 is the unique least Frobenius norm least squares solution of Problem 2.1.

Proof. Assume that the initial value $(\widetilde{A}(1), \widetilde{B}(1), \widetilde{C}(1))$ is chosen as in (2.23). According to Theorem 2.17, the solution ( $\left.\widetilde{A}_{*}, \widetilde{B}_{*}, \widetilde{C}_{*}\right)$ can be obtained by Algorithm 2.11 and it has the following form:

$$
\left\{\begin{align*}
\widetilde{A}_{*} & =Y\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Y^{H}-S_{n}\left(Y\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Y^{H}\right) S_{n} \\
& -E_{s} E_{s}^{T}\left[Y\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Y^{H}-S_{n}\left(Y\left(X \Lambda^{2}\right)^{H}+X \Lambda^{2} Y^{H}\right) S_{n}\right] E_{s} E_{s}^{T} \\
\widetilde{B}_{*} & =Y(X \Lambda)^{H}+X \Lambda Y^{H}-S_{n}\left(Y(X \Lambda)^{H}+X \Lambda Y^{H}(k)\right) S_{n}  \tag{2.24}\\
& -E_{t} E_{t}^{T}\left[Y(X \Lambda)^{H}+X \Lambda Y^{H}-S_{n}\left(Y(X \Lambda)^{H}+X \Lambda Y^{H}\right) S_{n}\right] E_{t} E_{t}^{T} \\
\widetilde{C}_{*} & =Y X^{H}+X Y^{H}-S_{n}\left(Y X^{H}+X Y^{H}\right) S_{n} \\
& -E_{u} E_{u}^{T}\left[Y X^{H}+X Y^{H}-S_{n}\left(Y X^{H}+X Y^{H}\right) S_{n}\right] E_{u} E_{u}^{T}
\end{align*}\right.
$$

Now suppose that $(\widetilde{A}, \widetilde{B}, \widetilde{C})$ is an arbitrary solution of Problem 2.1. It follows from Lemma 2.18 that there exists $\left(\mathcal{W}_{1}, \mathcal{W}_{2}, \mathcal{W}_{3}\right) \in \widetilde{S}_{1} \times \widetilde{S}_{2} \times \widetilde{S}_{3}$ such that

$$
\begin{equation*}
(\widetilde{A}, \widetilde{B}, \widetilde{C})=\left(\widetilde{A}_{*}+\mathcal{W}_{1}, \widetilde{B}_{*}+\mathcal{W}_{2}, \widetilde{C}_{*}+\mathcal{W}_{3}\right) \tag{2.25}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathcal{W}_{1} X \Lambda^{2}+\mathcal{W}_{2} X \Lambda+\mathcal{W}_{3} X=O \tag{2.26}
\end{equation*}
$$

By the relation (2.26) and Lemmas 2.6 and 2.7, we have

$$
\begin{aligned}
& \left\langle\widetilde{A}_{*}, \mathcal{W}_{1}\right\rangle+\left\langle\widetilde{B}_{*}, \mathcal{W}_{2}\right\rangle+\left\langle\widetilde{C}_{*}, \mathcal{W}_{3}\right\rangle=4\left\langle Y\left(X \Lambda^{2}\right)^{T}, \mathcal{W}_{1}\right\rangle+4\left\langle Y(X \Lambda)^{T}, \mathcal{W}_{2}\right\rangle+4\left\langle Y X^{T}, \mathcal{W}_{3}\right\rangle \\
= & 4\left\langle Y, \mathcal{W}_{1} X \Lambda^{2}+\mathcal{W}_{2} X \Lambda+\mathcal{W}_{3} X\right\rangle=0
\end{aligned}
$$

which, together with (2.25), yields

$$
\begin{aligned}
& \|\widetilde{A}\|^{2}+\|\widetilde{B}\|^{2}+\|\widetilde{C}\|^{2}=\left\|\widetilde{A}_{*}+\mathcal{W}_{1}\right\|^{2}+\left\|\widetilde{B}_{*}+\mathcal{W}_{2}\right\|^{2}+\left\|\widetilde{C}_{*}+\mathcal{W}_{3}\right\|^{2} \\
= & \left\|\widetilde{A}_{*}\right\|^{2}+\left\|\widetilde{B}_{*}\right\|^{2}+\left\|\widetilde{C}_{*}\right\|^{2}+\left\|\mathcal{W}_{1}\right\|^{2}+\left\|\mathcal{W}_{2}\right\|^{2}+\left\|\mathcal{W}_{3}\right\|^{2} \\
& +2\left[\left\langle\widetilde{A}_{*}, \mathcal{W}_{1}\right\rangle+\left\langle\widetilde{B}_{*}, \mathcal{W}_{2}\right\rangle+\left\langle\widetilde{C}_{*}, \mathcal{W}_{3}\right\rangle\right] \\
= & \left\|\widetilde{A}_{*}\right\|^{2}+\left\|\widetilde{B}_{*}\right\|^{2}+\left\|\widetilde{C}_{*}\right\|^{2}+\left\|\mathcal{W}_{1}\right\|^{2}+\left\|\mathcal{W}_{2}\right\|^{2}+\left\|\mathcal{W}_{3}\right\|^{2} \geq\left\|\widetilde{A}_{*}\right\|^{2}+\left\|\widetilde{B}^{*}\right\|^{2}+\left\|\widetilde{C}_{*}\right\|^{2} .
\end{aligned}
$$

This implies that the solution $\left(\widetilde{A}_{*}, \widetilde{B}_{*}, \widetilde{C}_{*}\right)$ is the least Frobenius norm solution of Problem 2.1. The proof is completed.

## 3 Iterative Method for Problem 1.2

In this section, we will establish an iterative method for solving Problem 1.2. For given $\bar{A}, \bar{B}, \bar{C} \in \mathbb{C}^{n \times n}$ and an arbitrary solution $(A, B, C)=\left(\widetilde{A}+\widetilde{A}_{p}, \widetilde{B}+\widetilde{B}_{q}, \widetilde{C}+\widetilde{C}_{r}\right) \in S_{E}$. Let $\bar{A}^{\prime}=\bar{A}-\widetilde{A}_{p}, \bar{B}^{\prime}=\bar{B}-\widetilde{B}_{q}$ and $\bar{C}^{\prime}=\bar{C}-\widetilde{C}_{r}$. Then

$$
\begin{align*}
& \|A-\bar{A}\|^{2}+\|B-\bar{B}\|^{2}+\|C-\bar{C}\|^{2}=\left\|\widetilde{A}-\bar{A}^{\prime}\right\|^{2}+\left\|\widetilde{B}-\bar{B}^{\prime}\right\|^{2}+\left\|\widetilde{C}-\bar{C}^{\prime}\right\|^{2} \\
& =\left\|\widetilde{A}-\frac{\bar{A}^{\prime}+\bar{A}^{\prime H}}{2}\right\|^{2}+\left\|\widetilde{B}-\frac{\bar{B}^{\prime}+\bar{B}^{\prime H}}{2}\right\|^{2}+\left\|\widetilde{C}-\frac{\bar{C}^{\prime}+\bar{C}^{\prime} H}{2}\right\|^{2}+\left\|\frac{\bar{A}^{\prime}-\bar{A}^{\prime H}}{2}\right\|^{2} \\
& +\left\|\frac{\bar{B}^{\prime}-\bar{B}^{\prime H}}{2}\right\|^{2}+\left\|\frac{\bar{C}^{\prime}-\bar{C}^{\prime} H}{2}\right\|^{2} \\
& =\left\|\widetilde{A}-\frac{\bar{A}^{\prime}+\bar{A}^{\prime} H-S_{n}\left(\bar{A}^{\prime}+\bar{A}^{\prime H}\right) S_{n}}{4}\right\|^{2}+\left\|\widetilde{B}-\frac{\bar{B}^{\prime}+\bar{B}^{\prime}{ }^{H}-S_{n}\left(\bar{B}^{\prime}+\bar{B}^{\prime}{ }^{H}\right) S_{n}}{4}\right\|^{2} \\
& +\left\|\widetilde{C}-\frac{\bar{C}^{\prime}+\bar{C}^{\prime} H}{}-S_{n}\left(\bar{C}^{\prime}+\bar{C}^{\prime H}\right) S_{n}\right\|^{2}+\left\|\frac{\bar{A}^{\prime}-\bar{A}^{\prime} H}{2}\right\|^{2} \\
& +\left\|\frac{\bar{B}^{\prime}-\bar{B}^{\prime} H}{2}\right\|^{2}+\left\|\frac{{\overline{C^{\prime}}}^{-}-\bar{C}^{\prime} H}{2}\right\|^{2} \\
& =\left\|\widetilde{A}-\frac{\bar{A}^{\prime}+\bar{A}^{\prime} H-S_{n}\left(\bar{A}^{\prime}+\bar{A}^{\prime} H\right) S_{n}-E_{s} E_{s}^{T}\left[\bar{A}^{\prime}+\bar{A}^{\prime H}-S_{n}\left(\bar{A}^{\prime}+\bar{A}^{\prime} H\right) S_{n}\right] E_{s} E_{s}^{T}}{4}\right\|^{2} \\
& +\left\|\widetilde{B}-\frac{\bar{B}^{\prime}+\bar{B}^{\prime}{ }^{H}-S_{n}\left(\bar{B}^{\prime}+\bar{B}^{\prime H}\right) S_{n}-E_{t} E_{t}^{T}\left[\bar{B}^{\prime}+\bar{B}^{\prime H}-S_{n}\left(\bar{B}^{\prime}+\bar{B}^{\prime H}\right) S_{n}\right] E_{t} E_{t}^{T}}{4}\right\|^{2} \\
& +\left\|\widetilde{C}-\frac{\bar{C}^{\prime}+\bar{C}^{\prime} H}{}-S_{n}\left(\bar{C}^{\prime}+\bar{C}^{\prime} H\right) S_{n}-E_{u} E_{u}^{T}\left[\bar{C}^{\prime}+\bar{C}^{\prime H}-S_{n}\left(\bar{C}^{\prime}+\bar{C}^{\prime H}\right) S_{n}\right] E_{u} E_{u}^{T}\right\|^{2} \\
& +\left\|\frac{E_{s} E_{s}^{T}\left[\bar{A}^{\prime}+\bar{A}^{\prime H}-S_{n}\left(\bar{A}^{\prime}+\bar{A}^{\prime H}\right) S_{n}\right] E_{s} E_{s}^{T}}{4}\right\|^{2} \\
& +\left\|\frac{E_{t} E_{t}^{T}\left[\bar{B}^{\prime}+\bar{B}^{\prime}{ }^{H}-S_{n}\left(\bar{B}^{\prime}+\bar{B}^{\prime}{ }^{H}\right) S_{n}\right] E_{t} E_{t}^{T}}{4}\right\|^{2} \\
& \left.+\| \frac{E_{u} E_{u}^{T}\left[\bar{C}^{\prime}+\bar{C}^{\prime} H\right.}{}-S_{n}\left(\bar{C}^{\prime}+\bar{C}^{\prime}{ }^{H}\right) S_{n}\right] E_{u} E_{u}^{T} \|^{2} \\
& +\left\|\frac{\bar{A}^{\prime}-\bar{A}^{\prime H}}{2}\right\|^{2}+\left\|\frac{\bar{B}^{\prime}-\bar{B}^{\prime H}}{2}\right\|^{2}+\left\|\frac{\bar{C}^{\prime}-\bar{C}^{\prime H}}{2}\right\|^{2} . \tag{3.1}
\end{align*}
$$

Let

$$
\begin{align*}
& \mathcal{A}=\frac{\bar{A}^{\prime}+\bar{A}^{\prime H}-S_{n}\left(\bar{A}^{\prime}+\bar{A}^{\prime} H\right) S_{n}-E_{s} E_{s}^{T}\left[\bar{A}^{\prime}+\bar{A}^{\prime H}-S_{n}\left(\bar{A}^{\prime}+\bar{A}^{\prime H}\right) S_{n}\right] E_{s} E_{s}^{T}}{4},  \tag{3.2}\\
& \mathcal{B}=\frac{\bar{B}^{\prime}+\bar{B}^{\prime H}-S_{n}\left(\bar{B}^{\prime}+\bar{B}^{\prime} H\right) S_{n}-E_{t} E_{t}^{T}\left[\bar{A}^{\prime}+\bar{B}^{\prime H}-S_{n}\left(\bar{B}^{\prime}+\bar{B}^{\prime} H\right) S_{n}\right] E_{t} E_{t}^{T}}{4}, \tag{3.3}
\end{align*}
$$

and

$$
\begin{equation*}
\mathcal{C}=\frac{\bar{C}^{\prime}+\bar{C}^{\prime H}-S_{n}\left(\bar{C}^{\prime}+\bar{C}^{\prime}\right) S_{n}-E_{u} E_{u}^{T}\left[\bar{A}^{\prime}+\bar{C}^{\prime H}-S_{n}\left(\bar{C}^{\prime}+\bar{C}^{\prime} H\right) S_{n}\right] E_{u} E_{u}^{T}}{4} . \tag{3.4}
\end{equation*}
$$

It is easy to see that $(\mathcal{A}, \mathcal{B}, \mathcal{C}) \in \widetilde{S}_{1} \times \widetilde{S}_{2} \times \widetilde{S}_{3}$. And, it follows from (3.1)-(3.4) that

$$
\left.\left.\begin{array}{rl} 
& \|A-\bar{A}\|^{2}+\|B-\bar{B}\|^{2}+\|C-\bar{C}\|^{2} \\
= & \left.\|\widetilde{A}-\mathcal{A}\|^{2}+\|\widetilde{B}-\mathcal{B}\|^{2}+\|\widetilde{C}-\mathcal{C}\|^{2}+\| \frac{E_{s} E_{s}^{T}\left[\bar{A}^{\prime}+\bar{A}^{\prime} H\right.}{}-S_{n}\left(\bar{A}^{\prime}+\bar{A}^{\prime} H\right) S_{n}\right] E_{s} E_{s}^{T} \\
& +\left\|\frac{E_{t} E_{t}^{T}\left[\bar{B}^{\prime}+\bar{B}^{\prime} H\right.}{}\right\|^{2}-S_{n}\left(\bar{B}^{\prime}+\bar{B}^{\prime} H\right. \\
4
\end{array}\right) S_{n}\right] E_{t} E_{t}^{T} \|^{2} .
$$

By direct calculations, we have

$$
\begin{aligned}
& (\widetilde{A}, \widetilde{B}, \widetilde{C}) \in \min _{\widetilde{S}_{1} \times \widetilde{S}_{2} \times \widetilde{S}_{3}}\left\|\widetilde{A} X \Lambda^{2}+\widetilde{B} X \Lambda+\widetilde{C} X-\widetilde{Z}\right\| \\
& =\min _{(\widetilde{A}, \widetilde{B}, \widetilde{C}) \in \widetilde{S}_{1} \times \widetilde{S}_{2} \times \widetilde{S}_{3}}\left\|(\widetilde{A}-\mathcal{A}) X \Lambda^{2}+(\widetilde{B}-\mathcal{B}) X \Lambda+(\widetilde{C}-\mathcal{C}) X-\left(\widetilde{Z}-\mathcal{A} X \Lambda^{2}-\mathcal{B} X \Lambda-\mathcal{C} X\right)\right\|
\end{aligned}
$$

Denote $\widetilde{\mathcal{Z}}=\widetilde{Z}-\mathcal{A} X \Lambda^{2}-\mathcal{B} X \Lambda-\mathcal{C} X$. Then

$$
\begin{align*}
& \min _{(\widetilde{A}, \widetilde{B}, \widetilde{C}) \in \widetilde{S}}^{1 \times \widetilde{S}_{2} \times \widetilde{S}_{3}} \\
&=\left\|\widetilde{A} X \Lambda^{2}+\widetilde{B} X \Lambda+\widetilde{C} X-\widetilde{Z}\right\|  \tag{3.6}\\
&(\widetilde{A}-\mathcal{A}, \widetilde{B}-\mathcal{B}, \widetilde{C}-\mathcal{C}) \in \widetilde{S}_{1} \times \widetilde{S}_{2} \times \widetilde{S}_{3}
\end{align*}\left\|(\widetilde{A}-\mathcal{A}) X \Lambda^{2}+(\widetilde{B}-\mathcal{B}) X \Lambda+(\widetilde{C}-\mathcal{C}) X-\widetilde{\mathcal{Z}}\right\| .
$$

According to (3.5) and (3.6), it follows that

$$
\min _{(A, B, C) \in S_{E}}\left[\|A-\bar{A}\|^{2}+\|B-\bar{B}\|^{2}+\|C-\bar{C}\|^{2}\right]
$$

is equivalent to finding a least Frobenius norm solution of the following problem

$$
\begin{equation*}
\min _{(\widetilde{\mathcal{A}}, \widetilde{\mathcal{B}}, \widetilde{\mathcal{C}}) \in \widetilde{S}_{1} \times \widetilde{S}_{2} \times \widetilde{S}_{3}}\left\|\widetilde{\mathcal{A}} X \Lambda^{2}+\widetilde{\mathcal{B}} X \Lambda+\widetilde{\mathcal{C}} X-\widetilde{\mathcal{Z}}\right\| \tag{3.7}
\end{equation*}
$$

where

$$
\begin{aligned}
& \widetilde{\mathcal{A}}=\widetilde{A}-\mathcal{A}=\widetilde{A}-\frac{\bar{A}^{\prime}+\bar{A}^{\prime} H-S_{n}\left(\bar{A}^{\prime}+\bar{A}^{\prime} H\right) S_{n}-E_{s} E_{s}^{T}\left[\bar{A}^{\prime}+\bar{A}^{\prime H}-S_{n}\left(\bar{A}^{\prime}+\bar{A}^{\prime} H\right) S_{n}\right] E_{s} E_{s}^{T}}{4} \\
& \left.\widetilde{\mathcal{B}}=\widetilde{B}-\mathcal{B}=\widetilde{B}-\frac{\bar{B}^{\prime}+\bar{B}^{\prime H}-S_{n}\left(\bar{B}^{\prime}+\bar{B}^{\prime} H\right) S_{n}-E_{t} E_{t}^{T}\left[\bar{B}^{\prime}+\bar{B}^{\prime} H\right.}{4}-S_{n}\left(\bar{B}^{\prime}+\bar{B}^{\prime} H\right) S_{n}\right] E_{t} E_{t}^{T} \\
& \widetilde{\mathcal{C}}=\widetilde{C}-\mathcal{C}=\widetilde{C}-\frac{\bar{C}^{\prime}+\bar{C}^{\prime} H}{}-S_{n}\left(\bar{C}^{\prime}+\bar{C}^{\prime H}\right) S_{n}-E_{u} E_{u}^{T}\left[\bar{C}^{\prime}+\bar{C}^{\prime H}-S_{n}\left(\bar{C}^{\prime}+\bar{C}^{\prime} H\right) S_{n}\right] E_{u} E_{u}^{T} \\
& 4
\end{aligned}
$$

By applying Algorithm 2.11, we can obtain the unique least Frobenius norm solution ( $\widetilde{\mathcal{A}}_{*}, \widetilde{\mathcal{B}}_{*}, \widetilde{\mathcal{C}}_{*}$ ) of the problem (3.7). Then the unique solution of Problem 1.2 can be obtained as follows:

$$
\left\{\begin{array}{l}
\widehat{A}=\widetilde{\mathcal{A}}_{*}+\frac{\bar{A}^{\prime}+\bar{A}^{\prime} H-S_{n}\left(\bar{A}^{\prime}+\bar{A}^{\prime} H\right) S_{n}-E_{s} E_{s}^{T}\left[\bar{A}^{\prime}+\bar{A}^{\prime H}-S_{n}\left(\bar{A}^{\prime}+\bar{A}^{\prime} H\right) S_{n}\right] E_{s} E_{s}^{T}}{4}+\widetilde{A}_{p} \\
\widehat{B}=\widetilde{\mathcal{B}}_{*}+\frac{\bar{B}^{\prime}+\bar{B}^{\prime H}-S_{n}\left(\bar{B}^{\prime}+\bar{B}^{\prime H}\right) S_{n}-E_{t} E_{t}^{T}\left[\bar{B}^{\prime}+\bar{B}^{\prime H}-S_{n}\left(\bar{B}^{\prime}+\bar{B}^{\prime H}\right) S_{n}\right] E_{t} E_{t}^{T}}{4}+\widetilde{B}_{q} \\
\widehat{C}=\widetilde{\mathcal{C}}_{*}+\frac{\bar{C}^{\prime}+\bar{C}^{\prime H}-S_{n}\left(\bar{C}^{\prime}+\bar{C}^{\prime H}\right) S_{n}-E_{u} E_{u}^{T}\left[\bar{C}^{\prime}+\bar{C}^{\prime H}-S_{n}\left(\bar{C}^{\prime}+\bar{C}^{\prime H}\right) S_{n}\right] E_{u} E_{u}^{T}}{4}+\widetilde{C}_{r}
\end{array}\right.
$$

## 4 Numerical Experiments

In this section, we report some numerical results to illustrate the efficiency of the proposed method and verify conclusions in this paper. All of the tests were run on the Intel (R) Core (TM), where the CPU is 2.40 GHz and the memory is 8.0 GB . The programming language is MATLAB R2015a. In view of the influence of round-off errors, we regard a matrix $T$ as the zero matrix if $\langle T, T\rangle<10^{-10}$.

Example 4.1. Consider the following quadratic inverse eigenvalue problem

$$
A X \Lambda^{2}+B X \Lambda+C X=O
$$

where

| $X=$ |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: |
| ( $0.8147+0.6948 i$ | $0.2785+0.7655 i$ | $0.9572+0.7094 i$ | $0.7922+0.1190 i$ | $0.6787+0.7513 i$ | $0.7060+0.5472 i$ |
| $0.9058+0.3171 i$ | $0.5469+0.7952 i$ | $0.4854+0.7547 i$ | $0.9595+0.4984 i$ | $0.7577+0.2551 i$ | $0.0318+0.1386 i$ |
| $0.1270+0.9502 i$ | $0.9575+0.1869 i$ | $0.8003+0.2760 i$ | $0.6557+0.9597 i$ | $0.7431+0.5060 i$ | $0.2769+0.1493 i$ |
| $0.9134+0.0344 i$ | $0.9649+0.4898 i$ | $0.1419+0.6797 i$ | $0.0357+0.3404 i$ | $0.3922+0.6991 i$ | $0.0462+0.2575 i$ |
| $0.6324+0.4387 i$ | $0.1576+0.4456 i$ | $0.4218+0.6551 i$ | $0.8491+0.5853 i$ | $0.6555+0.8909 i$ | $0.0971+0.8407 i$ |
| ( $0.0975+0.3816 i$ | $0.9706+0.6463 i$ | $0.9157+0.1626 i$ | $0.9340+0.2238 i$ | $0.1712+0.9593 i$ | $0.8235+0.2543 i$ |
| $\begin{aligned} & \text { and } \Lambda=\operatorname{diag}(4.6176+0.2816 i, 1.1012+0.4456 i, 0.4056+0.2760 i,-0.4480+0.1190 i, \\ & -0.0826+0.0344 i, 0.1136+0.2238 i) . \quad \text { Let } s=\{2,5\}, t=\{2,5\}, u=\{3,4\}, \quad A_{p}= \end{aligned}$ |  |  |  |  |  |
| $\begin{aligned} & \left(\begin{array}{cc} 0.25+0.01 i & -0.5+0.2 i \\ -0.25-0.01 i & 0.25+0.01 i \end{array}\right), \quad B_{q}=\left(\begin{array}{cc} 0.25 & -0.5+0.2 i \\ -0.25-0.01 i & 0.25+0.01 i \end{array}\right) \text { and } C_{r}= \\ & \left(\begin{array}{cc} 0.05+0.01 i & 0.15+0.03 i \\ 0.2+0.02 i & 0.1+0.01 i \end{array}\right) \text {. Then } \end{aligned}$ |  |  |  |  |  |
|  |  |  |  |  |  |
|  | $\widetilde{A}_{p}=\left(\begin{array}{l}0 \\ 0 \\ 0 \\ 0 \\ 0 \\ 0\end{array}\right.$ | 0 $0.25+0.01 i$ 0 0 $-0.25-0.01 i$ 0 | $\begin{array}{rrr}0 & 0 & 0 \\ 0 & 0 & -0.5+ \\ 0 & 0 & 0 \\ 0 & 0 & 0 \\ 0 & 0 & 0.25+ \\ 0 & 0 & 0\end{array}$ | +0.2i $\left.\begin{array}{ll}0 \\ & 0 \\ & 0 \\ 0.01 i & 0 \\ & 0\end{array}\right)$ |  |
|  | $\widetilde{B}_{q}=\left(\begin{array}{l}0 \\ 0 \\ 0 \\ 0 \\ 0 \\ 0\end{array}\right.$ | 0 0.25 0 0 $-0.25-0.01 i$ 0 | $\begin{array}{rrr}0 & 0 & 0 \\ 0 & 0 & -0.5+ \\ 0 & 0 & 0 \\ 0 & 0 & 0 \\ 0 & 0 & 0.25+0 \\ 0 & 0 & 0\end{array}$ | +0.2i $\left.\begin{array}{cc}0 \\ & 0 \\ & 0 \\ 0.01 i & 0 \\ & 0\end{array}\right)$ |  |
|  | $\widetilde{C}_{r}=\left(\begin{array}{l}0 \\ 0 \\ 0 \\ 0 \\ 0 \\ 0\end{array}\right.$ | $\begin{array}{cc}0 & 0 \\ 0 & 0 \\ 0 & 0.05+0.01 i \\ 0 & 0.2+0.02 i \\ 0 & 0 \\ 0 & 0\end{array}$ | $\begin{array}{cc} & 0 \\ & 0 \\ i & 0.15+0.03 i \\ 0.1+0.01 i \\ & 0 \\ & 0\end{array}$ | $\left.\begin{array}{lll}0 & 0 \\ 0 & 0 \\ 0 & 0 \\ 0 & 0 \\ 0 & 0 \\ 0 & 0\end{array}\right)$, |  |
| $\widetilde{Z}=-\widetilde{A}_{p} X \Lambda^{2}-\widetilde{B}_{q} X \Lambda-\widetilde{C}_{r} X=$ |  |  |  |  |  |
|  | - | 0 | 0 | 0 |  |
| ( $\begin{gathered}4.6417+0.6790 i \\ -0.1328-0.0813 i\end{gathered}$ | $0.0901+0.0309 i$ $-0.1760-0.1213 i$ | $0.0834+0.1384 i$ $-0.0381-0.1280 i$ | $-0.0790+0.0019 i$ $-0.0183-0.1067 i$ | $-0.0316-0.0103 i$ $-0.0700-0.1494 i$ | $-0.0847+0.0854 i$ $-0.0116-0.0502 i$ |
| -0.1328-0.08132 | -0.1760-0.12132 | $-0.0381-0.12802$ $-0.1619-0.1406 i$ | -0.1121-0.2395i | -0.1707-0.1899i | ${ }^{-0.0544-0.0616 i}$ |
| $\xrightarrow[0]{1.8788-0.5150 i}$ | $0.0682+0.3270 i$ | $-0.0054+0.0201 i$ | $-0.0072+0.0057 i$ | $\mathrm{ol}_{0}^{0.0021+0.0131 i}$ | ${ }_{0}^{0.0477-0.0160 i}$ |

First, we illustrate the efficiency of Algorithm 2.11 for solving the quadratic inverse eigenvalue problem. Let the initial matrices $\widetilde{A}(1)=\widetilde{B}(1)=\widetilde{C}(1)=\operatorname{zeros}(6,6)$. By Algorithm 2.11, we obtain the least Frobenius norm solution of Problem 1.1 after 69 iterative steps as follows.


At this moment, the norms of $R(k)$ and $P(k)$ are $\|R(k)\|=2.1513$ and $\|P(k)\|=7.4137 e-11$, respectively. The relationship between the number of iterations and the norm of $P(k)$ is shown in Figure 1. Figure 1 illustrates that Algorithm 2.11 is efficient for solving the constrained quadratic inverse eigenvalue problem.


Figure 1: The relationship between the number of iterations and the norm of $P(k)$ for Example 4.1

Second, we compare Algorithm 2.11 with the CGNR method developed by Hajarian and Abbas [24]. For fairness, we also choose the zero matrices as the initial values. By the CGNR method, we obtain the solution of Problem 2.1 after 67 iterative steps as follows.


In other words, the submatrix constrained solution generated by CGNR method can only be approximated.

By contrast, our proposed method projects the iterative sequences $\left\{\widetilde{A}_{\mathrm{NCG}}(k)\right\},\left\{\widetilde{B}_{\mathrm{NCG}}(k)\right\}$ and $\left\{\widetilde{C}_{\text {NCG }}(k)\right\}$ to linear subspaces $\widetilde{S}_{1}, \widetilde{S}_{2}$ and $\widetilde{S}_{3}$, respectively. Hence, we always have

$$
\widetilde{A}_{\mathrm{NCG}}(k)[s \mid s]=O, \widetilde{B}_{\mathrm{NCG}}(k)[t \mid t]=O, \widetilde{C}_{\mathrm{NCG}}(k)[u \mid u]=O
$$

This illustrates that our proposed method can solve the submatrix constrained quadratic inverse eigenvalue problem more accurately.

## 5 Concluding Remarks

This paper is concerned with the least squares solution of a class of constrained quadratic inverse eigenvalue problem and its optimal approximation problem. We propose a nonlinear conjugate gradient method for finding the solution over Hermitian centroskew matrices with a submatrix constraint. The convergence analysis of the proposed method is given. Numerical results illustrate that the proposed method is efficient for the quadratic inverse eigenvalue problem.

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