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On the perturbation and subproper splittings for the generalized inverse $A_{T,S}^{(2)}$ of rectangular matrix $A^{\stackrel{t}{\sim}}$

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Abstract

In this paper, the perturbation and subproper splittings for the generalized inverse $A_{T,S}^{(2)}$, the unique matrix X such that XAX = X, R(X) = T and N(X) = S, are considered. We present lower and upper bounds for the perturbation of $A_{T,S}^{(2)}$. Convergence of subproper splittings for computing the special solution $A_{T,S}^{(2)}b$ of restricted rectangular linear system Ax = b, $x \in T$, are studied. For the solution $A_{T,S}^{(2)}b$ we develop a characterization. Therefore, we give a unified treatment of the related problems considered in literature by Ben-Israel, Berman, Hanke, Neumann, Plemmons, etc. © 2001 Published by Elsevier Science B.V.

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1. Introduction and preliminary results

In this paper, we adopt the same notations on generalized inverses used in Ben-Israel and Greville [3] and Berman and Plemmons [6]. It is well known that the commonly important six kinds of generalized inverses: the Moore–Penrose inverse A^+ , the weighted Moore–Penrose inverse A^+_{MN} , the Drazin inverse A^D , the group inverse A_g , the Bott–Duffin inverse $A_{(L)}^{(-1)}$ and the generalized

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Bott–Duffin inverse $A_{(L)}^{(+)}$ can all be viewed as special {2}-inverses $A_{T,S}^{(2)}$, having the prescribed range T and null space S.

The {2}-inverse has many applications, for example, the application in the iterative methods for solving nonlinear equations [3,18] and the applications to statistics [11,13]. In particular, {2}-inverse plays an important role in stable approximations of ill-posed problems and in linear and nonlinear problems involving rank-deficient generalized inverse [17,27]. There are many numerical methods for computing $A_{T,S}^{(2)}$ (see [8,30] and references therein). In 1966, Ben-Israel [1] gave the perturbation for the Moore–Penrose inverse A^+ , and recently

In 1966, Ben-Israel [1] gave the perturbation for the Moore–Penrose inverse A^+ , and recently this result has been extended to the Drazin inverse A^D and the group inverse A_g , see [28,31]. In this paper, our first goal is to make a unified treatment for the perturbation of the generalized inverse $A_{rs}^{(2)}$.

In the past decades, many authors studied various kinds of linear equations. Ben-Israel [2], Verghese [24], and Wang [25] considered the Cramer rule for minimum norm solution or least-squares solutions of consistent and inconsistent linear equations

$$Ax = b$$
,

where $A \in \mathbb{C}^{m \times n}$, and $b \in \mathbb{C}^m$. Werner [33] presented the Cramer rule for restricted linear equations:

$$Ax = b, \quad x \in K,$$

where $A \in \mathbb{C}^{m \times n}$, $b \in \mathbb{C}^m$, and K is a complementary subspace of N(A). Wang [26] discussed the special singular equations

$$Ax = b, \quad x \in R(A^k),$$

where $A \in \mathbb{C}^{n \times n}$, k = index(A), and $b \in R(A^k)$. Summarizing these equations, Chen [7,9] discussed a general restricted linear equation:

$$Ax = b, \quad x \in T, \tag{1.1}$$

where $A \in \mathbb{C}^{m \times n}$, and T is a subspace of \mathbb{C}^n . In [7] Chen showed that the restricted system (1.1) has a solution if and only if $b \in AT$, and has a unique solution if and only if $b \in AT$ and $T \cap N(A) = \{0\}$. Let subspace $S \in \mathbb{C}^m$ be such that $AT \oplus S = \mathbb{C}^m$, then one sees that $A_{T,S}^{(2)}b$ is a solution or the unique solution of system (1.1) assuming it is consistent. Our second goal is to establish a characterization of this special solution despite system (1.1) is solvable or not.

To obtain a solution of (1.1) and its special cases, one often splits A into the form

$$A = U - V. \tag{1.2}$$

Then there results the natural iteration

$$x_{i+1} = U^- V x_i + U^- b, (1.3)$$

where U^- is some generalized inverse of *A*. In literature, various choices of *U* are considered in order to compute the corresponding solution of (1.1), see [4,5,14,15,19,29] and Section 4. Our last contribution to this problem is to consider two kinds of subproper splittings (1.2) of *A* for computing the solution $A_{T,S}^{(2)}b$ of (1.1), which, to our best knowledge, include and extend all kinds of existing splittings for solution of (1.1).

The following lemmata are needed in what follows.

Lemma 1.1 (Ben-Israel and Greville [3]). Let $A \in \mathbb{C}^{m \times n}$, rank(A) = r, and let T and S be subspaces of \mathbb{C}^n and \mathbb{C}^m , respectively, dim $T = \dim S^{\perp} = t \leq r$. Then A has a $\{2\}$ -inverse X such that R(X) = Tand N(X) = S if and only if

$$AT \oplus S = \mathbb{C}^m, \tag{1.4}$$

in which case X is unique and denoted by $X = A_{T,S}^{(2)}$.

Note 1 (Ben-Israel and Greville [3]). When t = r in Lemma 1.1, condition (1.4) is equivalent to

$$R(A) \oplus S = \mathbb{C}^m \quad \text{and} \quad T \oplus N(A) = \mathbb{C}^n,$$
(1.5)

in which case $A_{T,S}^{(2)}$ is indeed $A_{T,S}^{(1,2)}$, the unique {1,2}-inverse of A having range T and null space S.

Note 2 (Ben-Israel and Greville [3]). $AA_{T,S}^{(2)} = P_{AT,S}, A_{T,S}^{(2)}A = P_{T,(A^*S^{\perp})^{\perp}}, AA_{T,S}^{(1,2)} = P_{R(A),S}$ and $A_{T,S}^{(1,2)}A = P_{R(A),S}$ $P_{T,N(A)}$.

Lemma 1.2 (Chen [8]). (1) Let $A \in \mathbb{C}^{m \times n}$. Then, for the Moore–Penrose inverse A^+ and the weighted Moore–Penrose inverse A_{MN}^+ , one has

- (i) A⁺ = A^(1,2)_{R(A^{*}),N(A^{*})} = A⁽²⁾_{R(A^{*}),N(A^{*})}.
 (ii) A⁺_{M,N} = A^(1,2)_{R(A[#]),N(A[#])} = A⁽²⁾_{R(A[#]),N(A[#])}, where A[#] = N⁻¹A^{*}M. N and M are Hermitian positive-definite matrices of order n and m, respectively.
- (2) Let $A \in \mathbb{C}^{n \times n}$. Then, for the Drazin inverse A^{D} , the group inverse A_{g} , the Bott–Duffin inverse $A_{(L)}^{(-1)}$ and the generalized Bott–Duffin inverse $A_{(L)}^{(+)}$, one has
- (iii) $A^{D} = A^{(2)}_{R(A^{k}), N(A^{k})}$ where k = index(A).
- (iv) In particular, when index(A) = 1, $A_g = A_{R(A),N(A)}^{(2)}$.
- (v) $A_{(L)}^{(-1)} = A_{L,L^{\perp}}^{(2)} = P_L(AP_L + P_{L^{\perp}})^{-1}$, where *L* is a subspace of \mathbb{C}^n such that $AL \oplus L^{\perp} = \mathbb{C}^n$ and P_L is the orthogonal projector on *L*; (vi) $A_{(L)}^{(+)} = A_{(S)}^{(-1)} = A_{S,S^{\perp}}^{(2)}$ where $S = R(P_LA)$.

2. Perturbation theory

Let $B = A + E \in \mathbb{C}^{m \times n}$. Let subspaces $T, \tilde{T} \subseteq \mathbb{C}^n$ and subspaces $S, \tilde{S} \in \mathbb{C}^m$ be such that $AT \oplus S = \mathbb{C}^m$ and $B\tilde{T} \oplus \tilde{S} = \mathbb{C}^m$. To present perturbation bounds for generalized inverse $A_{T,S}^{(2)}$ we give first a simple expression of $B_{\tilde{T},\tilde{S}}^{(2)}$ in terms of $A_{T,S}^{(2)}$ and E.

Theorem 2.1. Let $B = A + E \in \mathbb{C}^{m \times n}$. Let subspaces T, \tilde{T} of \mathbb{C}^n and subspaces S, \tilde{S} of \mathbb{C}^m be such that $AT \oplus S = \mathbb{C}^m$ and $B\tilde{T} \oplus \tilde{S} = \mathbb{C}^m$. Then

$$B_{\tilde{T},\tilde{S}}^{(2)} = (I + A_{T,S}^{(2)}E)^{-1}A_{T,S}^{(2)} = A_{T,S}^{(2)}(I + EA_{T,S}^{(2)})^{-1}$$
(2.1)

if and only if

$$\tilde{T} = T \quad and \quad \tilde{S} = S.$$
 (2.2)

Proof. The necessary part is obvious. Assume (2.2) holds. We first show that $I + A_{T,S}^{(2)}E$ is invertible. If not, then there is a nonzero vector x such that

$$(I + A_{TS}^{(2)}E)x = 0$$

i.e.,

$$x = -A_{T,S}^{(2)} Ex \in T = \tilde{T}.$$
(2.3)

Multiplying $A_{T,S}^{(2)}A$ with (2.3) yields

 $A_{Ts}^{(2)}Bx = 0,$

which implies $Bx \in S = \tilde{S}$. But from (2.3) we see $Bx \in B\tilde{T}$. Hence Bx = 0. Therefore, $x = B_{\tilde{\tau},\tilde{S}}^{(2)}Bx = 0$ by $x \in \tilde{T}$ and Note 2, a contradiction.

By direct computation, we have

$$(I + A_{T,S}^{(2)}E)B_{T,S}^{(2)} = B_{T,S}^{(2)} + A_{T,S}^{(2)}BB_{T,S}^{(2)} - A_{T,S}^{(2)}AB_{T,S}^{(2)}$$
$$= B_{T,S}^{(2)} + A_{T,S}^{(2)} - B_{T,S}^{(2)}$$
$$= A_{T,S}^{(2)}.$$

Thus $B_{T,S}^{(2)} = (I + A_{T,S}^{(2)}E)^{-1}A_{T,S}^{(2)}$. Since $(I + A_{T,S}^{(2)}E)A_{T,S}^{(2)} = A_{T,S}^{(2)}(I + EA_{T,S}^{(2)})$, so $B_{T,S}^{(2)} = A_{T,S}^{(2)}(I + EA_{T,S}^{(2)})^{-1}$ and the proof is complete.

Corollary 2.2 (Wei [32]). Let $B=A+E \in \mathbb{C}^{m \times n}$, and let T, S be subspaces of \mathbb{C}^n and \mathbb{C}^m , respectively, such that $AT \oplus S = \mathbb{C}^m$. Suppose $R(E) \subseteq AT$ and $R(E^*) \subseteq A^*S^{\perp}$. If $||EA_{T,S}^{(2)}|| < 1$, then $B_{T,S}^{(2)}$ exists and

$$B_{T,S}^{(2)} = (I + A_{T,S}^{(2)}E)^{-1}A_{T,S}^{(2)} = A_{T,S}^{(2)}(I + EA_{T,S}^{(2)})^{-1}.$$

Proof. The hypothesis $||EA_{T,S}^{(2)}|| < 1$ implies that $I + A_{T,S}^{(2)}E$ and $I + EA_{T,S}^{(2)}$ are nonsingular. By the assumption and Note 2,

$$B = A(I + A_{T,S}^{(2)}E) = (I + EA_{T,S}^{(2)})A$$
(2.4)

holds and $BT \subseteq AT$. Let P_T be the orthogonal projector on T, then $BT = R(BP_T)$ and $AT = R(AP_T)$. It follows from (2.4) that

$$BP_T = (I + EA_{T,S}^{(2)})AP_T.$$

Thus $\operatorname{rank}(BP_T) = \operatorname{rank}(AP_T)$ and BT = AT.

Notice that $AT \oplus S = \mathbb{C}^m$, then $BT \oplus S = \mathbb{C}^m$ and $B_{T,S}^{(2)}$ exist. We conclude the proof of the corollary applying Theorem 2.1. \Box

From Theorem 2.1 and Corollary 2.2, we immediately obtain the following results.

Corollary 2.3. Let B = A + E. Let subspaces $T \subseteq \mathbb{C}^n$ and $S \subseteq \mathbb{C}^m$ be such that $T \oplus N(A) = \mathbb{C}^n$ and $R(A) \oplus S = \mathbb{C}^m$. Suppose $R(E) \subseteq R(A)$ and $R(E^*) \subseteq R(A^*)$. If $||EA_{T,S}^{(1,2)}|| < 1$, then $B_{T,S}^{(1,2)}$ exists and $B_{T,S}^{(1,2)} = (I + A_{T,S}^{(1,2)}E)^{-1}A_{T,S}^{(1,2)} = A_{T,S}^{(1,2)}(I + EA_{T,S}^{(1,2)})^{-1}.$ (2.5) **Corollary 2.4** (Ben-Israel [1]). Let B = A + E. Suppose $R(E) \subseteq R(A)$ and $R(E^*) \subseteq R(A^*)$. If $||EA^+||$ < 1, then

$$B^{+} = (I + A^{+}E)^{-1}A^{+} = A^{+}(I + EA^{+})^{-1}.$$
(2.6)

Corollary 2.5 (Wei [29]). Let $B = A + E \in \mathbb{C}^{n \times n}$ with index(A) = k. Suppose $R(E) \subseteq R(A^k)$ and $R(E^*) \subseteq R(A^{k^*})$. If $||EA^{D}|| < 1$, then index(B) = k and $B^{\rm D} = (I + A^{\rm D}E)^{-1}A^{\rm D} = A^{\rm D}(I + EA^{\rm D})^{-1}$ (2.7)

We are in a position to present the main result of this section.

Theorem 2.6. Let $B = A + E \in \mathbb{C}^{m \times n}$. Let subspaces $T \subseteq \mathbb{C}^n$ and $S \subseteq \mathbb{C}^m$ be such that $AT \oplus S = \mathbb{C}^m$. Suppose $R(E) \subseteq AT$ and $R(E^*) \subseteq A^*S^{\perp}$. If $\triangle = ||EA_{T,S}^{(2)}|| < 1$, then

$$\frac{\triangle}{(1+\triangle)\mathscr{K}(A)} \leqslant \frac{\|B_{T,S}^{(2)} - A_{T,S}^{(2)}\|}{\|A_{T,S}^{(2)}\|} \leqslant \frac{\triangle}{1-\triangle},\tag{2.8}$$

where $\mathscr{K}(A) = ||A|| ||A_{TS}^{(2)}||$ is the condition number of generalized inverse $A_{TS}^{(2)}$.

Proof. From Theorem 2.1 and Corollary 2.2 we have

$$B_{T,S}^{(2)} - A_{T,S}^{(2)} = -B_{T,S}^{(2)} E A_{T,S}^{(2)}.$$
(2.9)

Taking norm of (2.9) leads to

$$\|B_{T,S}^{(2)} - A_{T,S}^{(2)}\| \le \Delta \|B_{T,S}^{(2)}\|.$$
(2.10)

It follows from (2.1) that

$$\frac{\|A_{T,S}^{(2)}\|}{1+\Delta} \leqslant \|B_{T,S}^{(2)}\| \leqslant \frac{\|A_{T,S}^{(2)}\|}{1-\Delta}.$$
(2.11)

Substituting (2.11) into (2.10), one has

$$\frac{\|B_{T,S}^{(2)} - A_{T,S}^{(2)}\|}{\|A_{T,S}^{(2)}\|} \leqslant \frac{\triangle}{1 - \triangle}.$$
(2.12)

On the other hand, we note that $BB_{T,S}^{(2)} = P_{BT,S} = P_{AT,S} = AA_{T,S}^{(2)}$ since AT = BT, so we have

$$EA_{T,S}^{(2)} = B(A_{T,S}^{(2)} - B_{T,S}^{(2)}).$$
(2.13)

Substituting (2.4) into (2.13) we have

$$EA_{T,S}^{(2)} = (I + EA_{T,S}^{(2)})A(A_{T,S}^{(2)} - B_{T,S}^{(2)}).$$
(2.14)

It follows from (2.14) that

$$\frac{\|B_{T,S}^{(2)} - A_{T,S}^{(2)}\|}{\|A_{T,S}^{(2)}\|} \ge \frac{\triangle}{(1+\triangle)\mathscr{K}(A)}$$
(2.15)

and, considering (2.12) and (2.15), the proof is complete. \Box

Corollary 2.7 (Ben-Israel [1]). Let $B = A + E \in \mathbb{C}^{m \times n}$. Suppose $R(E) \subseteq R(A)$ and $R(E^*) \subseteq R(A^*)$. If $\Delta = ||EA^+|| < 1$, then

$$\frac{\triangle}{(1+\triangle)\mathscr{K}(A)} \leqslant \frac{\|B^+ - A^+\|}{\|A^+\|} \leqslant \frac{\triangle}{1-\triangle},\tag{2.16}$$

where $\mathscr{K}(A) = ||A|| ||A^+||$ is the condition number of Moore–Penrose inverse A^+ .

Corollary 2.8 (Wei [29]). Let $B = A + E \in \mathbb{C}^{n \times n}$ with index(A) = k. Suppose $R(E) \subseteq R(A^k)$ and $R(E^*) \subseteq R(A^{k^*})$. If $\Delta = ||EA^D|| < 1$, then

$$\frac{\triangle}{(1+\triangle)\mathscr{K}(A)} \leqslant \frac{\|B^{\mathrm{D}} - A^{\mathrm{D}}\|}{\|A^{\mathrm{D}}\|} \leqslant \frac{\triangle}{1-\triangle},\tag{2.17}$$

where $\mathscr{K}(A) = ||A|| ||A^{D}||$ is the condition number of Drazin inverse A^{D} .

3. Characterization of $A_{T,S}^{(2)}b$

In this section, we shall establish a characteristic for the special solution $A_{T,S}^{(2)}b$ of the restricted linear equations (1.1) despite $b \in AT$ or not. Our result gives better characterization of $A_{T,S}^{(2)}b$ than that of Chen [7], and extends the results of Hanke and Neumann [12] for $A_{T,S}^{(1,2)}b$ and that of Ben-Israel [3] for A^+b .

Recall that $AA_{T,S}^{(2)} = P_{AT,S}$ and $A_{T,S}^{(2)}A = P_{T,(A^*S^{\perp})^{\perp}}$. As done in [10,12], we define the norms

$$\|x\|_{AT,S}^{2} = \|P_{AT,S}x\|_{2}^{2} + \|(I - P_{AT,S})x\|_{2}^{2} \quad \forall x \in \mathbb{C}^{m}$$
(3.1)

and

$$\|y\|_{T,(A^*S^{\perp})^{\perp}}^2 = \|P_{T,(A^*S^{\perp})^{\perp}}y\|_2^2 + \|(I - P_{T,(A^*S^{\perp})^{\perp}})y\|_2^2 \quad \forall y \in \mathbb{C}^n.$$
(3.2)

Theorem 3.1. The solution $A_{T,S}^{(2)}b$ is a least-squares solution of Ax = b in N(A) + T with respect to the norm $\|.\|_{AT,S}$, i.e.,

$$\|b - AA_{T,S}^{(2)}b\|_{AT,S} = \min_{x \in N(A)+T} \|b - Ax\|_{AT,S}.$$
(3.3)

Moreover, $A_{T,S}^{(2)}b$ is the unique minimum $\|.\|_{T,(A^*S^{\perp})^{\perp}}$ -norm solution to $A_{T,S}^{(2)}Ax = A_{T,S}^{(2)}b$.

Proof. For any $x \in N(A) + T$, we have

$$\|b - Ax\|_{AT,S}^{2} = \|P_{AT,S}(b - Ax)\|_{2}^{2} + \|(I - P_{AT,S})(b - Ax)\|_{2}^{2}$$

$$= \|P_{AT,S}b - Ax\|_{2}^{2} + \|(I - P_{AT,S})b\|_{2}^{2}$$

$$\ge \|(I - P_{AT,S})b\|_{2}^{2} = \|b - AA_{T,S}^{(2)}b\|_{AT,S}^{2}, \qquad (3.4)$$

which implies $A_{T,S}^{(2)}b$ is a least-squares solution Au = b in N(A) + T with respect to the norm $\|.\|_{AT,S}$.

It is evident that $A_{T,S}^{(2)}Ax = A_{T,S}^{(2)}b$ is consistent and the general solution is $x = A_{T,S}^{(2)}b + z$, $\forall z \in (A^*S^{\perp})^{\perp}$. Thus

$$\|x\|_{T,(A^*S^{\perp})^{\perp}}^2 = \|P_{T,(A^*S^{\perp})^{\perp}}(A_{T,S}^{(2)}b+z)\|_2^2 + \|(I-P_{T,(A^*S^{\perp})^{\perp}})(A_{T,S}^{(2)}b+z)\|_2^2$$

$$= \|A_{T,S}^{(2)}b\|_2^2 + \|z\|_2^2$$

$$\ge \|A_{T,S}^{(2)}b\|_2^2 = \|A_{T,S}^{(2)}b\|_{T,(A^*S^{\perp})^{\perp}}^2.$$
(3.5)

The equality holds in (3.5) if and only if z=0 which implies $A_{T,S}^{(2)}b$ is the minimum $\|.\|_{T,(A^*S^{\perp})^{\perp}}$ -norm solution of $A_{T,S}^{(2)}Ax = A_{T,S}^{(2)}b$. \Box

Corollary 3.2 (Hanke and Neumann [12]). Let $A \in \mathbb{C}^{m \times n}$, and let subspaces $T \subseteq \mathbb{C}^n$ and $S \subseteq \mathbb{C}^m$ be such that $T \oplus N(A) = \mathbb{C}^n$ and $R(A) \oplus S = \mathbb{C}^m$. Then $A_{T,S}^{(1,2)}b$ is a least-squares solution of Ax = b with respect to the norm $\|.\|_{R(A),S}$; moreover, $A_{T,S}^{(1,2)}b$ is the unique minimum $\|.\|_{T,N(A)}$ -norm solution of $A_{T,S}^{(1,2)}Ax = A_{T,S}^{(1,2)}b$.

For $T = R(A^*)$ and $S = N(A^*)$, we obtain

Corollary 3.3 (Ben-Israel and Greville [3]). Let $A \in \mathbb{C}^{m \times n}$. Then A^+b is the minimum norm least-squares solution of Ax = b.

4. Subproper splittings

Iterative methods for solution of system (1.1) of the form

$$x_{i+1} = Bx_i + c, \quad i = 0, 1, 2, \dots,$$
 (4.1)

where *B* is the *n*th-order complex matrix, are often employed. For this reason *B* is commonly called the iteration matrix. As mentioned in Section 1, *B* arises in a large number of cases from a splitting (1.2) of *A*.

The purpose of this section is to unify and extend well-known results concerning the convergence of iterative scheme (4.1) to a special solution of (1.1), that is, the solution $A_{T,S}^{(2)}b$.

In our study, we find it instructive to classify the known results into the following categories:

- (i) m = n, A and U in (1.2) are nonsingular, and the iteration matrix of (4.1) is $B = U^{-1}V$ (see [23,34,20,16,22]).
- (ii) m = n, in (1.2) A is a singular matrix and U is a nonsingular matrix and the iteration matrix is $B = U^{-1}V$ (see [15,21]).
- (iii) A and U in (1.2) are rectangular matrices and R(U) = R(A) and N(U) = N(A). This splitting is called proper splitting by Berman and Plemmons. In this case, the iteration matrix is $B = U^+V$ (see [4,5]).
- (iv) A and U in (1.2) are rectangular matrices, $m \ge n$, U has full column rank, and $UU^+A = A$, and the iteration matrix is $B = U^+V$ (see [14]).

- (v) m = n, A and U in (1.2) satisfy $R(U) = R(A^k)$ and $N(U) = N(A^k)$ where k = index(A). This splitting is called index splitting by Wei. In this case the iteration matrix is $U_a V$ (see [29]).
- (vi) A and U in (1.2) are rectangular matrices. Neumann call splitting (1.2) is subproper if $R(A) \subseteq R(U)$ and $R(A^*) \subseteq R(U^*)$. In this case the iteration matrix of (4.1) is $B = U^+V$ (see [19]).

To unify all kinds of splittings above we give

Definition 4.1. Let $A=U-V \in \mathbb{C}^{m \times n}$. Let subspaces $T, \tilde{T} \subseteq \mathbb{C}^n$ and $S, \tilde{S} \subseteq \mathbb{C}^m$ be such that $AT \oplus S = \mathbb{C}^m$ and $U\tilde{T} \oplus \tilde{S} = \mathbb{C}^m$. The splitting (1.2) is called subproper if $T \subseteq \tilde{T}$ and $\tilde{S} \subseteq S$, and is called proper if $T = \tilde{T}$ and $S = \tilde{S}$.

It is easy to see that the splittings from (i)-(vi) but (v) are all subproper splittings according to Definition 4.1. We shall remark that the concept of subproper splitting in Definition 4.1 is more general than that of Neumann.

Now we give

Definition 4.2. Let $A, U \in \mathbb{C}^{m \times n}$. Let subspaces $T, \tilde{T} \subseteq \mathbb{C}^n$ and $S, \tilde{S} \subseteq \mathbb{C}^m$ be such that $AT \oplus S = \mathbb{C}^m$ and $U\tilde{T} \oplus \tilde{S} = \mathbb{C}^m$. Splitting (1.2) is called subproper if $\tilde{T} \subseteq T$ and $S \subseteq \tilde{S}$, and is called proper if $T = \tilde{T}$ and $S = \tilde{S}$.

The index splitting (v) is a subproper splitting in Definition 4.2. To make a distinction between the subproper splittings in Definitions 4.1 and 4.2, we shall call them in the sequel types I and II, respectively. In both cases the iterative scheme (4.1) becomes

$$x_{i+1} = U_{\tilde{T},\tilde{S}}^{(2)} V x_i + U_{\tilde{T},\tilde{S}}^{(2)} b, \quad i = 0, 1, 2, \dots$$
(4.2)

Theorem 4.3. Let $A = U - V \in \mathbb{C}^{m \times n}$, and let $T, \tilde{T} \subseteq \mathbb{C}^n$ and $S, \tilde{S} \subseteq \mathbb{C}^m$ be such that $AT \oplus S = \mathbb{C}^m$ and $U\tilde{T} \oplus S = \mathbb{C}^m$. Let splitting (1.2) be of type I or type II subproper splitting. Then (i)

$$A_{T,S}^{(2)} = (I - U_{\tilde{T},\tilde{S}}^{(2)}V)^{-1}U_{\tilde{T},\tilde{S}}^{(2)} = U_{\tilde{T},\tilde{S}}^{(2)}(I - VU_{\tilde{T},\tilde{S}}^{(2)})^{-1}$$
(4.3)

if and only if $T = \tilde{T}$ and $S = \tilde{S}$, i.e., the splitting is proper.

(ii) The sequence of the iterations of (4.2) converges to $A_{T,S}^{(2)}b$ for every $x_0 \in \mathbb{C}^n$ and for every $b \in \mathbb{C}^m$ if and only if $\rho(U_{\tilde{\tau}\tilde{S}}^{(2)}V) < 1$ and splitting (1.2) is proper.

Proof. Part (i) follows directly from Theorem 2.1, we have only to prove (ii). Assume that $\rho(U_{\tilde{T},\tilde{S}}^{(2)}V)$ < 1 and $T = \tilde{T}$, $S = \tilde{S}$. Then by (i)

$$(I - U_{\tilde{T},\tilde{S}}^{(2)}V)^{-1}U_{\tilde{T},\tilde{S}}^{(2)} = A_{T,S}^{(2)}.$$

From (4.2) it is easily proven by induction that

$$x_{i} = (U_{\tilde{T},\tilde{S}}^{(2)}V)^{i}x_{0} + \sum_{j=0}^{i-1} (U_{\tilde{T},\tilde{S}}^{(2)}V)^{j}U_{\tilde{T},\tilde{S}}^{(2)}b.$$

$$(4.4)$$

Since $\rho(U_{\tilde{\tau},\tilde{s}}^{(2)}V) < 1$, it follows from (4.4) that

$$\lim_{i\to\infty} x_i = (I - U_{\tilde{T},\tilde{S}}^{(2)}V)^{-1}U_{\tilde{T},\tilde{S}}^{(2)}b = A_{T,S}^{(2)}b.$$

Conversely, assume that the sequence of $\{x_i\}$ with respect to (4.2) converges to $A_{T,S}^{(2)}b$ independent of the initial guess x_0 and b, we must have $\rho(U_{\tilde{T},\tilde{S}}^{(2)}V) < 1$ and

$$A_{T,S}^{(2)} = (I - U_{\tilde{T},\tilde{S}}^{(2)}V)^{-1}U_{\tilde{T},\tilde{S}}^{(2)} = U_{\tilde{T},\tilde{S}}^{(2)}(I - VU_{\tilde{T},\tilde{S}}^{(2)})^{-1}$$

By Theorem 2.1 we have $T = \tilde{T}$ and $S = \tilde{S}$ and the proof is complete. \Box

Corollary 4.4 (Berman and Plemmons [5]). Let $A = U - V \in \mathbb{C}^{m \times n}$ be such that R(U) = R(A) and N(U) = N(A). Then

- (i) $A^+ = (I U^+ V)^{-1} U^+ = U^+ (I V U^+)^{-1}$
- (ii) The sequence of iterations

$$x_{i+1} = U^+ V x_i + U^+ b$$

converges to A^+b for every $x_0 \in \mathbb{C}^n$ and every $b \in \mathbb{C}^m$ if and only if $\rho(U^+V) < 1$.

Corollary 4.5 (Wei [29]). Let $A \in \mathbb{C}^{n \times n}$ with index(A) = k. Let A = U - V be such that $R(U) = R(A^k)$ and $N(U) = N(A^k)$. Then

- (i) $A^{\mathrm{D}} = (I U_g V)^{-1} U_g = U_g (I V U_g)^{-1}.$
- (ii) The sequence of iterations

$$x_{i+1} = U_g V x_i + U_g b$$

converges to $A^{\mathrm{D}}b$ for every $x_0 \in \mathbb{C}^n$ and every $b \in \mathbb{C}^n$ if and only if $\rho(U_q V) < 1$.

To present convergence theorem for type I subproper splitting, we assume that $b \in AT$ for the moment.

Lemma 4.6. Let subspaces $T, \tilde{T} \subseteq \mathbb{C}^n$ and $S, \tilde{S} \subseteq \mathbb{C}^m$ be such that $AT \oplus S = \mathbb{C}^m$ and $U\tilde{T} \oplus \tilde{S} = \mathbb{C}^m$. Let $b \in AT$ and splitting (1.2) be type I subproper splitting of A. Then every solution $x \in \tilde{T}$ (of which at least one exists, e.g., $A_{T,S}^{(2)}b$) satisfies $x = U_{\tilde{T},\tilde{S}}^{(2)}Vx + U_{\tilde{T},\tilde{S}}^{(2)}b$.

Proof. The thesis follows immediately on applying $U_{\tilde{T},\tilde{S}}^{(2)}$ to Ux = Vx + b. \Box

Let $x \in \tilde{T}$ be any fixed solution of Ax = b and define the "errors"

$$e_i = x_i - x. \tag{4.5}$$

Then from (4.2) and Lemma 4.6, the recursion satisfied by these errors is

$$e_i = Be_{i-1} = \dots = B^i e_0, \tag{4.6}$$

where $B = U_{\tilde{T},\tilde{S}}^{(2)}V$ is the iteration matrix. From this relation we immediately have the following theorem.

Theorem 4.7. Under the conditions of Lemma 4.6, the following three statements are equivalent:

- (i) For any x_0 the sequence $\{x_i\}$ of (4.2) converges to a solution of Ax = b in \tilde{T} .
- (ii) For every e_0 the error sequence $\{e_i\}$ of (4.6) converges to a vector in $N(A) \cap \tilde{T}$. (iii) The iteration matrix $B = U_{\tilde{T},\tilde{S}}^{(2)}V$ is semi-convergent, i.e.,

$$\rho(B) \leqslant 1, \tag{4.7}$$

$$\lambda \in \sigma(B) \quad and \quad |\lambda| = 1 \Rightarrow \lambda = 1,$$
(4.8)

$$\operatorname{index}(I-B) \leqslant 1. \tag{4.9}$$

Proof. The equivalence of (i) and (ii) and (iii) are evident. \Box

When splitting (1.2) is a type I subproper splitting and the iteration matrix B is semi-convergent, then iteration (4.2) will converge to

$$(I-B)_g U_{\tilde{T},\tilde{S}}^{(2)} b + [I-(I-B)(I-B)_g] x_0.$$
(4.10)

For the convergence of type II subproper splitting, we choose the initial vector $x_0 \in \tilde{T}$, then the iterative scheme (4.2) is equivalent to

$$x_{i+1} = x_i + U_{\tilde{T},\tilde{S}}^{(2)}(b - Ax_i)$$
(4.11)

for $i = 0, 1, 2, \dots$.

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We now introduce a relaxation parameter β and consider

$$x_{i+1} = x_i + \beta U_{\tilde{T},\tilde{S}}^{(2)}(b - Ax_i).$$
(4.12)

Theorem 4.8. Let $A=U-V \in \mathbb{C}^{m \times n}$. Let subspaces $T, \tilde{T} \subseteq \mathbb{C}^n$ and $S, \tilde{S} \subseteq \mathbb{C}^m$ be such that $AT \oplus S = \mathbb{C}^m$ and $U\tilde{T} \oplus \tilde{S} = \mathbb{C}^m$. Let splitting (1.2) be a type II subproper splitting of A and $\beta > 0$. Then

(i) For any initial guess $x_0 \in T$, the sequence $\{x_i\}$ of (4.12) converges to $A_{T,S}^{(2)}b$ if and only if

$$\rho(P_T - \beta U_{\tilde{T},\tilde{S}}^{(2)}A) < 1.$$
(4.13)

(ii) Let (4.13) holds and $x_i \to x$ as $i \to \infty$. Then $x = A_{T,S}^{(2)}b$ has the representation

$$x = \beta (P_{T^{\perp}} + \beta U_{\tilde{T},\tilde{S}}^{(2)} A)^{-1} U_{\tilde{T},\tilde{S}}^{(2)} b.$$
(4.14)

Moreover, if $||P_T - \beta U_{\tilde{T},\tilde{S}}^{(2)}A|| = q < 1$, then x_i has the error estimation

$$\|x_i - A_{T,S}^{(2)}b\| \leq q^i (\|x_0\| + \beta (1-q)^{-1} \|U_{\tilde{T},\tilde{S}}^{(2)}b\|).$$
(4.15)

(iii) If all eigenvalues of $U_{\tilde{\tau},\tilde{s}}^{(2)}A$ are nonnegative, then (4.13) holds if and only if

$$\operatorname{rank}(U^{(2)}_{\tilde{T},\tilde{S}}A) = \dim T \quad and \quad 0 < \beta < \frac{2}{\lambda_{\max}(U^{(2)}_{\tilde{T},\tilde{S}}A)},\tag{4.16}$$

where $\lambda_{\max}(U_{\tilde{T},\tilde{S}}^{(2)}A)$ is the largest eigenvalue of $U_{\tilde{T},\tilde{S}}^{(2)}A$.

Proof. Denote $A_{T,S}^{(2)}$ by G for simplicity. It is easy to observe that x = Gb satisfies $x = x + \beta U_{T,S}^{(2)}(b - Ax)$. Hence from (4.12) we have

$$\begin{aligned} x_i - Gb &= x_{i-1} - Gb - \beta U_{\tilde{T},\tilde{S}}^{(2)} A(x_{i-1} - Gb) \\ &= (P_T - \beta U_{\tilde{T},\tilde{S}}^{(2)} A)(x_{i-1} - Gb) \quad (\text{by } x_i \in T). \end{aligned}$$

Using this relation repeatedly, it follows that

$$x_i - Gb = (P_T - \beta U_{\tilde{T},\tilde{S}}^{(2)}A)^i (x_0 - Gb).$$
(4.17)

(i) If (4.13) holds, then from (4.17) we have obviously $x_i \to Gb$ as $i \to \infty$ for any $x_0 \in T$. Conversely, suppose that $x_i \to Gb$ for every $x_0 \in T$. This x_0 is expressible as $x_0 = P_T z$, z arbitrary, and now (4.17) can be written as

$$x_{i} - Gb = (P_{T} - \beta U_{\tilde{T},\tilde{S}}^{(2)}A)^{i}P_{T}(z - Gb)$$

= $(P_{T} - \beta U_{\tilde{T},\tilde{S}}^{(2)}AP_{T})^{i}(z - Gb).$ (4.18)

Let the unitary matrix $Q = [Q_1, Q_2]$ be such that $T = R(Q_1)$ and $T^{\perp} = R(Q_2)$. Then $P_T = Q_1 Q_1^*$. Therefore,

$$Q^{*}(P_{T} - \beta U_{\tilde{T},\tilde{S}}^{(2)}A)Q = \begin{bmatrix} I - \beta Q_{1}^{*}U_{\tilde{T},\tilde{S}}^{(2)}AQ_{1} & *\\ 0 & 0 \end{bmatrix},$$
(4.19)

$$Q^{*}(P_{T} - \beta U_{\tilde{T},\tilde{S}}^{(2)}AP_{T})Q = \begin{bmatrix} I - \beta Q_{1}^{*}U_{\tilde{T},\tilde{S}}^{(2)}AQ_{1} \ 0 \\ 0 \ 0 \end{bmatrix}$$
(4.20)

and

$$Q^{*}(x_{i} - Gb) = (Q^{*}(P_{T} - \beta U_{\tilde{T},\tilde{S}}^{(2)}A)Q)^{i}Q^{*}(z - Gb)$$

=
$$\begin{bmatrix} (I - \beta Q_{1}^{*}U_{\tilde{T},\tilde{S}}^{(2)}AQ_{1})^{i}\tilde{z} \\ 0 \end{bmatrix}, \quad \tilde{z} \text{ arbitrary.}$$
(4.21)

Since $x_i \to Gb$ as $i \to \infty$, it follows from (4.21) that $\rho(I - \beta U_{\tilde{I},\tilde{S}}^{(2)}A) < 1$. In view of (4.19), we have $\rho(P_T - \beta U_{\tilde{T},\tilde{S}}^{(2)}A) < 1.$ (ii) Note that

$$x_{i} = (P_{T} - \beta U_{\tilde{T},\tilde{S}}^{(2)}A)x_{i-1} + \beta U_{\tilde{T},\tilde{S}}^{(2)}b$$

$$= [I + (P_{T} - \beta U_{\tilde{T},\tilde{S}}^{(2)}A) + (P_{T} - \beta U_{\tilde{T},\tilde{S}}^{(2)}A)^{2} + \dots + (P_{T} - \beta U_{\tilde{T},\tilde{S}}^{(2)}A)^{i-1}]\beta U_{\tilde{T},\tilde{S}}^{(2)}b$$

$$+ (P_{T} - \beta U_{\tilde{T},\tilde{S}}^{(2)}A)^{i}x_{0}.$$
 (4.22)

If $\rho(P_T - \beta U_{\tilde{T},\tilde{S}}^{(2)}A) < 1$, then one has

$$x_{i} = \beta [I - (P_{T} - \beta U_{\hat{T}, \hat{S}}^{(2)} A)]^{-1} U_{\hat{T}, \hat{S}}^{(2)} b = \beta (P_{T^{\perp}} + \beta U_{\hat{T}, \hat{S}}^{(2)} A)^{-1} U_{\hat{T}, \hat{S}}^{(2)} b.$$

Furthermore, by (4.22) and (4.14), we get

$$x_i - Gb = -(P_T - \beta U_{\tilde{T},\tilde{S}}^{(2)}A)^i [I + (P_T - \beta U_{\tilde{T},\tilde{S}}^{(2)}A) + (P_T - \beta U_{\tilde{T},\tilde{S}}^{(2)}A)^2 + \cdots] \beta U_{\tilde{T},\tilde{S}}^{(2)}b$$

$$+(P_T-\beta U^{(2)}_{\tilde{T},\tilde{S}}A)^i x_0$$

giving (4.15) if $||P_T - \beta U_{\tilde{T},\tilde{S}}^{(2)}A|| = q < 1.$ (iii) For the proof of (iii) see [8, Theorem 2.1]. \Box

We should note that Chen proved an analogous result for the iterative scheme similar to (4.12)in the case of $b \in AT$ in [9]. However, in Theorem 4.7, we do not assume that (1.1) is consistent.

The following convergence theorem for type II subproper splitting is the direct result of Theorem 4.8.

Theorem 4.9. Let $A = U - V \in \mathbb{C}^{m \times n}$. Suppose subspaces $T, \tilde{T} \subseteq \mathbb{C}^n$ and $S, \tilde{S} \subseteq \mathbb{C}^m$ be such that $AT \oplus S = \mathbb{C}^m$ and $U\tilde{T} \oplus \tilde{S} = \mathbb{C}^m$. Let (1.2) be a type II subproper splitting of A. Then (i) If $\rho(P_T - U_{\tilde{T},\tilde{S}}^{(2)}A) < 1$, then the sequence $\{x_i\}$ of (4.2) converges to $A_{T,S}^{(2)}b$ provided $x_0 \in \tilde{T}$.

Moreover, if $||P_T - U_{\tilde{T},\tilde{S}}^{(2)}A|| = q < 1$, then x_i has the error estimation

$$|x_i - A_{T,S}^{(2)}b|| \leq q^i (||x_0|| + (1-q)^{-1} ||U_{\tilde{T},\tilde{S}}^{(2)}b||).$$
(4.23)

(ii) Suppose all the eigenvalues of $U_{\tilde{T},\tilde{S}}^{(2)}A$ are nonnegative and $0 < \lambda_{\max}(U_{\tilde{T},\tilde{S}}^{(2)}A) < 2$. If splitting (1.2) is proper, then $\rho(P_T - U_{\tilde{T},\tilde{S}}^{(2)}A) < 1$ and thus the sequence of iterations of (4.2) converges to $A_{TS}^{(2)}b$ for any $x_0 \in T$.

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