

OVERESTIMATION IN LINEAR INTERVAL EQUATIONS\*

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**Abstract.** Bounds are given for the overestimation of the solution set of a system of linear interval equations by various inclusion intervals for this set. In particular, a theorem of Gay [this Journal, 19 (1982), pp. 858-870] on the quadratic approximation property of the preconditioned fixed-point inverse is strengthened.

**Key words.** linear interval equations, overestimation, quadratic approximation

**AMS(MOS) subject classifications.** 65G10, 65F10

**1. Introduction.** A well-studied problem in interval analysis is the construction of enclosures for the set of solutions of the equations

$$(1) \quad \tilde{A}\tilde{x} = \tilde{b} \quad (\tilde{A} \in A, \tilde{b} \in b),$$

where  $A$  is a regular interval matrix (i.e. all  $\tilde{A} \in A$  are nonsingular) and  $b$  is an interval vector. While preconditioned Gauss elimination [Hansen [6]] and preconditioned interval iteration [Krawczyk [7]] turned out to be practically useful methods, the quality of the computed enclosures was paid little attention, apart from the special case of  $M$ -matrices. There, Barth and Nuding [2] showed that the limit  $A^{\infty}b$  of the Jacobi iteration for (1) agrees with the hull  $A^{\infty}b$  of the solution set of (1), and Barth and Nuding [2] and Beck [3] showed the same for the result  $A^G b$  of Gauss elimination provided that the right-hand side  $b$  satisfies  $0 \leq b, 0 \leq b$ , or  $0 \in b$ .

In the general case an asymptotic investigation of the overestimation of preconditioned Gauss elimination has been made by Müller [10] (cf. the exposition in Alefeld and Herzberger [1, Chap. 16]). Moreover, exponential worst case overestimation of Gauss elimination without preconditioning was demonstrated by Wongwises [15] and Schätzle [14] under natural simplifying assumptions. For the overestimation in case of preconditioned interval iteration a neat theorem of Gay [5] is available; however, his proof contains a gap. The present paper partly owes its existence to an attempt to repair the proof; in fact we prove a strengthened version of Gay's theorem (Theorem 3). We also prove a new optimality result for the midpoint inverse as preconditioning matrix (Theorem 4) and derive another enclosure for the hull  $A^{\infty}b$  with small overestimation, based on an approximate inverse  $C$  and an approximate solution  $\tilde{x}$  (Theorem 1). Using the concept of a skew vector, we are even able to approximate the solution set of (1) itself by a skew vector with small overestimation (Theorem 2).

The terminology follows Neumaier [12]; it is assumed that the reader is acquainted with §§ 1-4 and 7 of [12]. However, for convenience of the reader we repeat some basic notation. We denote by  $\mathbb{R}, \mathbb{R}^n$ , and  $\mathbb{R}^{n \times n}$  the set of real numbers,  $n$ -dimensional column vectors and  $n \times n$ -matrices, and by  $\mathbb{R}, \mathbb{R}^n$ , and  $\mathbb{R}^{n \times n}$  the set of real closed intervals,  $n$ -dimensional interval vectors and  $n \times n$  interval matrices, respectively. If  $x = [\underline{x}, \bar{x}] \in \mathbb{R}^n$  then we write  $\text{mid } x := (\underline{x} + \bar{x})/2$  for the midpoint,  $\text{rad } x := \rho(x) := (\bar{x} - \underline{x})/2$  for the radius,  $|\underline{x}| := \sup \{|\bar{x}|, |\underline{x}|\}$  for the absolute value, and  $\tilde{x} := x - \tilde{x} = [-\rho(x), \rho(x)]$  for the radial part of  $x$ ; similar definitions apply to interval

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<sup>1</sup> To emphasize the nature of  $C$  as a fixed real matrix, we do not write  $\tilde{C}$  which would be the "generic" notation.

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matrices. The Ostrowski operator  $(\cdot)$  associates with a square matrix  $A \in \mathbb{R}^{n \times n}$  the matrix  $A' = (A)$  with entries  $A'_{ik} := \min \{|\bar{a}|, |\underline{a}|\} \tilde{a}_{ik} := -|A_{ik}|$  for  $i \neq k$ . In terms of the Jacobi splitting  $A = D - E$  of  $A$ , defined by  $D_{ii} := A_{ii}, E_{ii} := 0$  and  $D_{ik} := 0, E_{ik} := -A_{ik}$  for  $i \neq k$  we have

$$(A) = (D) - |E|, \quad |A| = |D| + |E| \cong (D) + |E|.$$

Sum and difference of two bounded subsets  $\Sigma_1, \Sigma_2 \subseteq \mathbb{R}^n$  are defined as  $\Sigma_1 \pm \Sigma_2 := \{\tilde{x} \pm \tilde{y} \mid \tilde{x} \in \Sigma_1, \tilde{y} \in \Sigma_2\}$ , and the hull of a bounded subset  $\Sigma \subseteq \mathbb{R}^n$  is defined as  $\text{Hull } \Sigma := [\inf \Sigma, \sup \Sigma]$ , the smallest interval vector containing  $\Sigma$ . The (vector valued) distance of two interval vectors  $x, y \in \mathbb{R}^n$  is defined as

$$(2) \quad q(x, y) := \inf \{q \in \mathbb{R}^n \mid q \geq 0, x \subseteq y + [-q, q], y \subseteq x + [-q, q]\}.$$

One immediately establishes the following properties of the distance.

LEMMA 1. Let  $x, y, z \in \mathbb{R}^n$ . Then

$$\begin{aligned} q(x, y) = 0 & \text{ iff } x = y, \\ q(y, x) & = q(x, y), \\ q(x, z) & \leq q(x, y) + q(y, z). \end{aligned}$$

LEMMA 2. Let  $x, y \in \mathbb{R}^n$ . Then

$$\begin{aligned} (3) \quad x \subseteq y & \text{ iff } |\bar{x} - \bar{y}| \leq \rho(y) - \rho(x), \\ (4) \quad q(x, y) & = |\bar{x} - \bar{y}| + |\rho(x) - \rho(y)|. \end{aligned}$$

*Proof.* Equation (3) is well known, and the equivalence of (2) and (4) is proved in Lemma 1 of Neumaier [13].  $\square$

Note that often (4) is taken as a definition of the distance of intervals.

LEMMA 3. If  $f(\xi_1, \dots, \xi_m)$  is an arithmetic expression in  $m$  variables such that the variables  $\xi_1, \dots, \xi_m$  (where  $n \leq m$ ) occur only once then for  $x_1, \dots, x_m \in \mathbb{R}, \rho(x_i) = 0$  for  $i > n$ ,

$$f(x_1, \dots, x_m) = \{f(\xi_1, \dots, \xi_m) \mid \xi_i \in x_i (i = 1, \dots, m)\}.$$

In particular, if  $a, b \in \mathbb{R}^n$  then

$$a + b = \{\bar{a} + \bar{b} \mid \bar{a} \in a, \bar{b} \in b\},$$

and if  $A \in \mathbb{R}^{n \times n}, \tilde{x} \in \mathbb{R}^n$  then

$$A\tilde{x} = \{\tilde{A}\tilde{x} \mid \tilde{A} \in A\}.$$

*Proof.* Well known; see [1] or [11].  $\square$

We define a skew vector as a set of the form

$$(5) \quad \{\tilde{x} + C\tilde{r} \mid \tilde{r} \in r\},$$

where  $\tilde{x} \in \mathbb{R}^n, C \in \mathbb{R}^{n \times n}$  and  $\tilde{r} \in \mathbb{R}^n$ . Geometrically, a skew vector is (possibly degenerate) parallelepiped with axes parallel to the columns of  $C$ . Skew vectors occur in the literature on interval differential equations, hidden behind the concept of a moving coordinate frame (see e.g. Moore [11], Eijgenraam [4]). The hull of the skew vector (5) is easily seen to be the interval  $\tilde{x} + Cr$ . In case of large off-diagonal entries in  $C$  the distance between a skew vector and its hull may be quite large; cf. Fig. 1, which shows a skew vector (5) and its hull, where

$$\tilde{x} = \begin{pmatrix} 2 \\ 2 \end{pmatrix}, \quad C = \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}, \quad r = \begin{pmatrix} [-1, 1] \\ [-1, 1] \end{pmatrix}.$$

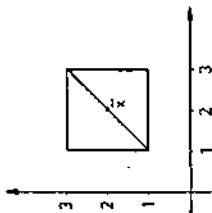


FIG. 1

Therefore it can be expected that, in general, the overestimation of a set  $\Sigma$  by an enclosing interval (e.g. its hull) can be considerably reduced by enclosing  $\Sigma$  in a suitable skew vector.

2. The solution set of linear interval equations. Let  $A \in \mathbb{R}^{n \times n}$ ,  $b \in \mathbb{R}^n$ , and suppose that  $A$  is regular. We are interested in narrow enclosures of the set

$$\Sigma(A, b) := \{\tilde{A}^{-1}\tilde{b} \mid \tilde{A} \in A, \tilde{b} \in b\}$$

of solutions of the linear system of interval equations  $\tilde{A}\tilde{x} = \tilde{b}$  ( $\tilde{A} \in A, \tilde{b} \in b$ ).

The optimal interval enclosure is

$$A''b := \Omega[\tilde{A}^{-1}\tilde{b} \mid \tilde{A} \in A, \tilde{b} \in b];$$

but in practice often only cruder enclosures are available. In this section we consider a particular enclosure based on the residual  $r := b - A\tilde{x}$  of an approximation  $\tilde{x}$  to the solution.

THEOREM 1. Suppose that  $C \in \mathbb{R}^{n \times n}$  is such that  $CA$  is regular. Let  $\tilde{x} \in \mathbb{R}^n$ ,  $r := b - A\tilde{x}$ , and suppose that  $d \in \mathbb{R}^n$  satisfies

$$A''((I - AC)r) \subseteq d \quad \text{or} \quad (CA)''((I - CA)(Cr)) \subseteq d.$$

Then  $x := \tilde{x} + Cr + d$  satisfies

$$(1) \quad A''b \subseteq x \subseteq A''b + (d - d);$$

in particular,

$$(2) \quad \rho(A''b, x) \leq 2\rho(d),$$

$$(3) \quad 0 \leq \text{rad } x - \text{rad } A''b \leq 2\rho(d).$$

Proof. By assumption,  $C$  and  $A$  are regular. Let  $\tilde{A} \in A$ ,  $\tilde{b} \in b$  and put

$$(4) \quad \tilde{r} := \tilde{b} - \tilde{A}\tilde{x}, \quad \tilde{y} := \tilde{x} + Cr, \quad \tilde{z} := \tilde{A}^{-1}\tilde{b}.$$

Then  $\tilde{z} - \tilde{y} = \tilde{A}^{-1}\tilde{b} - \tilde{x} - C(\tilde{b} - \tilde{A}\tilde{x}) = \tilde{A}^{-1}(I - \tilde{A}C)(\tilde{b} - \tilde{A}\tilde{x})$ , hence

$$(5) \quad \tilde{z} - \tilde{y} = \tilde{A}^{-1}(I - \tilde{A}C)\tilde{r} = (C\tilde{A})^{-1}(I - C\tilde{A})C\tilde{r} \in d.$$

Since  $\tilde{y} \in \tilde{x} + Cr$ , this implies (1). To prove (2) and (3), we use the fact that Lemma 3 implies

$$r = \{\tilde{b} - \tilde{A}\tilde{x} \mid \tilde{A} \in A, \tilde{b} \in b\};$$

hence every  $\tilde{r} \in r$  and every  $\tilde{y} \in$

$$(6) \quad \tilde{\Sigma} := \{\tilde{x} + C\tilde{r} \mid \tilde{r} \in r\}$$

can be written in the form (4). Therefore  $\Sigma(A, B) \subseteq \tilde{\Sigma} + d \subseteq \tilde{x} + Cr + d = x$ , and  $\tilde{\Sigma} \subseteq \Sigma(A, B) - d \subseteq A''b - d$ , whence  $\tilde{x} + Cr \subseteq A''b - d$ , and  $x = \tilde{x} + Cr + d \subseteq A''b - d + d$ . This implies (1), and (2). (3) are immediate consequences.  $\square$

For practical application,  $C$  should be an approximation of the inverse of some  $\tilde{A} \in A$  so that  $CA = I$ , and  $\tilde{x}$  should be an approximation of  $\tilde{A}^{-1}\tilde{b}$  for some  $\tilde{b} \in b$  so that  $A\tilde{x} = \tilde{b}$ . If the residuals  $I - CA$  and  $b - A\tilde{x}$  are of order  $O(\epsilon)$  then  $d$  can be taken of order  $O(\epsilon^2)$ , so that (1) overestimates the hull by only  $O(\epsilon^2)$ . Note that the amount of overestimation can be computed a posteriori from (3).

The proof of Theorem 1 suggests that  $\Sigma(A, b)$  itself may be enclosed with little overestimation by a skew vector. This is indeed the case.

THEOREM 2. Suppose that  $C \in \mathbb{R}^{n \times n}$  is such that  $AC$  is regular. Let  $\tilde{x} \in \mathbb{R}^n$ ,  $r := b - A\tilde{x}$  and suppose that

$$(AC)''((I - AC)r) \subseteq e \in \mathbb{R}^n.$$

Then  $\Sigma_0 := \{\tilde{x} + C\tilde{r} \mid \tilde{r} \in r + e\}$  satisfies

$$(7) \quad \Sigma(A, b) \subseteq \Sigma_0 \subseteq \Sigma(A, b) + |C|(e - e).$$

Proof. As before,

$$\tilde{e} := C^{-1}(\tilde{z} - \tilde{y}) = (\tilde{A}C)^{-1}(I - \tilde{A}C)\tilde{r} \in e,$$

so that  $\tilde{z} = \tilde{y} + C\tilde{e} = \tilde{x} + C(\tilde{r} + \tilde{e})$ , which implies  $\Sigma(A, b) \subseteq \Sigma_0$ . Again as before,  $\tilde{\Sigma} \subseteq \Sigma(A, b) - \{C\tilde{e} \mid \tilde{e} \in e\}$  so that  $\Sigma_0 \subseteq \tilde{\Sigma} + \{C\tilde{e} \mid \tilde{e} \in e\} \subseteq \Sigma(A, b) + \{C(\tilde{e} - \tilde{e}) \mid \tilde{e} \in e\} \subseteq \Sigma(A, b) + |C|(e - e)$ .  $\square$

Again, if the residuals  $I - AC$  and  $b - A\tilde{x}$  are of order  $O(\epsilon)$  then  $e$  can be chosen of order  $O(\epsilon^2)$ , so that (6) overestimates the solution set by only  $O(\epsilon^2)$ . Thus (6) can be expected to be a very accurate enclosure for the solution set. Note that the amount of overestimation can be calculated a posteriori from (7).

3. The preconditioned fixpoint inverse.  $A \in \mathbb{R}^{n \times n}$  is called an  $H$ -matrix if  $(A)_{ii} > 0$  for some  $u > 0$ ; cf. Neumaier [12, § 2.5]. In [12], the fixpoint inverse of an  $H$ -matrix  $A \in \mathbb{R}^{n \times n}$  was defined as the sublinear map  $A'' : \mathbb{R}^n \rightarrow \mathbb{R}^n$  which maps  $b \in \mathbb{R}^n$  to the unique solution  $x = A''b$  of the fixpoint equation

$$x = L^r(b + Ex),$$

where  $A = L - E$  is a direct splitting of  $A$ . By [12, Prop. 11], the vector  $x = A''b$  is an enclosure of the hull  $A''b$ . Here we consider the preconditioned equations

$$CA\tilde{x} = C\tilde{b} \quad (\tilde{A} \in A, \tilde{b} \in b),$$

and, assuming that  $CA$  is an  $H$ -matrix, the corresponding enclosure  $x = (CA)''(C\tilde{b})$  for  $(CA)''(C\tilde{b})$ . The norm used in the following is the maximum norm.

THEOREM 3. Suppose that  $C \in \mathbb{R}^{n \times n}$  is such that  $CA$  is an  $H$ -matrix. Then

$$(1) \quad z = A''b \subseteq (CA)''(C\tilde{b}) \subseteq (CA)^r(C\tilde{b}) \subseteq \tilde{z} + (CA)^{-1}|CA|\tilde{z}.$$

Moreover, if  $CA$  is strictly diagonally dominant, i.e.,

$$\beta := \max_i \frac{|CA|_{ii}}{\sum_{k \neq i} |CA|_{ki}} < 1$$

then

$$(2) \quad \|\text{rad } A''b\| \leq \|\text{rad } (CA)''(C\tilde{b})\| \leq \|\text{rad } (CA)^r(C\tilde{b})\| \leq \frac{1+\beta}{1-\beta} \|\text{rad } A''b\|.$$

*Proof.* Let  $CA = D - E$  be the Jacobi splitting of  $CA$ . Then  $x := (CA)^F(Cb)$  satisfies the equation

$$(3) \quad x = D^F(Cb + Ex).$$

Fix  $i \in \{1, \dots, n\}$ . The splitting equation implies

$$D_{ii} = \sum_j C_{ij}A_{jk}, \quad E_{ik} = -\sum_j C_{ij}A_{jk} \quad (k \neq i),$$

and (3) gives

$$x_i = \sum_k \frac{C_{ik}b_k + E_{ik}x_k}{D_{ii}}.$$

In these expressions, each variable apart from the  $C_{ij} \in \mathbb{R}$  occurs only once (note that  $i$  is fixed); hence these expressions equal the range (Lemma 3), and for each  $\xi \in x_i$  we can find  $\tilde{x} \in x$ ,  $\tilde{b} \in b$ ,  $\tilde{D} \in D$ ,  $\tilde{E} \in E$ ,  $\tilde{A} \in A$  such that  $CA = \tilde{D} - \tilde{E}$  and  $\tilde{y} := \tilde{D}^{-1}(Cb + \tilde{E}\tilde{x})$  satisfies  $\tilde{y}_i = \xi$ .

Now  $\tilde{z} := \tilde{A}^{-1}\tilde{b} \in A^{-1}b = z$  and  $\tilde{y} = \tilde{D}^{-1}(C\tilde{A}\tilde{z} + \tilde{E}\tilde{x}) = \tilde{D}^{-1}((\tilde{D} - \tilde{E})\tilde{z} + \tilde{E}\tilde{x}) = \tilde{z} + \tilde{D}^{-1}\tilde{E}(\tilde{x} - \tilde{z}) \in z + D^F E(x - z)$ ; therefore  $\xi \in z + D^F E(x - z)$ , and since  $\xi \in x_i$  and  $i$  were arbitrary we find

$$x \subseteq z + D^F E(x - z).$$

This implies

$$\begin{aligned} |x - \tilde{z}| &\subseteq |z + D^F E(x - z)| \subseteq \rho(z) + \langle D \rangle^{-1} |E| |x - z|, \\ \langle D \rangle |x - \tilde{z}| &\subseteq \langle D \rangle \rho(z) + |E| (|x - \tilde{z}| + \rho(z)). \end{aligned}$$

Therefore

$$(4) \quad |x - \tilde{z}| \subseteq (\langle D \rangle - |E|)^{-1} (\langle D \rangle + |E|) \rho(z).$$

In particular,

$$|x - \tilde{z}| \subseteq \langle CA \rangle^{-1} |CA| \rho(z), \quad x \in \tilde{z} + \langle CA \rangle^{-1} |CA| \xi$$

and (1) follows. Finally if  $CA$  is strictly diagonally dominant then  $B := \langle D \rangle^{-1} |E|$  satisfies  $\|B\| \leq \beta < 1$  whence

$$\begin{aligned} \| \text{rad } x \| &\leq \| x - \tilde{z} \| \leq \| (\langle D \rangle - |E|)^{-1} (\langle D \rangle + |E|) \rho(z) \| \\ &= \| (I - B)^{-1} (I + B) \rho(z) \| \leq \frac{1 + \beta}{1 - \beta} \| \rho(z) \| \end{aligned}$$

which implies (2).  $\square$

*Remarks.* 1) If  $\rho(A) = O(\epsilon)$  and  $C = \tilde{A}^{-1}$  then  $CA = I + \tilde{A}^{-1}\tilde{A} = I + O(\epsilon)$  whence  $\beta = O(\epsilon)$ . From (2) we therefore get

$$\frac{\text{rad}((CA)^F(Cb))}{\text{rad } A^{-1}b} = 1 + O(\epsilon),$$

i.e., very little overestimation.

2) For a computable bound on the overestimation we simply replace  $A^{-1}b$  in (2) by the slightly bigger  $(CA)^F(Cb)$ .

3) Theorem 3 implies that preconditioning with a good approximation  $C$  of  $\tilde{A}^{-1}$  (resulting in a small  $\beta$ , say  $\beta \leq \frac{1}{2}$ ) increases the radius of  $A^{-1}b$  only by a small factor  $\frac{1}{1 - \beta}$  (usually gives excellent enclosures. An asymptotic result of this type is already known for some time: Miller [10] proved

$$\text{rad}((CA)^F(Cb)) = \text{rad } A^{-1}b + O(\epsilon^2)$$

if  $C = \tilde{A}^{-1}$  and  $\rho(A)$ ,  $\rho(b) = O(\epsilon)$ ; cf. the exposition in Alefeld and Herzberger [1, Chap. 16].

4) Since  $q(x, z) = |\tilde{x} - \tilde{z}| + |\rho(x) - \rho(z)| = |x - \tilde{z}| - \rho(z)$  (note  $\rho(x) \geq \rho(z)$  since  $z \subseteq x$ ), we get from (4) the relation

$$q(x, z) \leq 2(\langle D \rangle - |E|)^{-1} |E| \rho(z),$$

which in the strictly diagonally dominant case implies

$$\|q((CA)^F(Cb), A^{-1}b)\| \leq \frac{2\beta}{1 - \beta} \|\text{rad } A^{-1}b\|.$$

5) If  $\|I - CA\| \leq \beta_0 < 1$  then the iteration  $x^{i+1} = Cb + Ex^i$ , where  $E := I - CA$ , converges to a limit  $x^*$  satisfying  $x^* = Cb + Ex^*$  and

$$z \subseteq (CA)^F(Cb) \subseteq x^*;$$

see Mayer [9], Gay [5] and Neumaier [12]. Proceeding as in the proof of Theorem 3 but with the splitting  $CA = I - E$  one obtains the relation

$$z \subseteq x^* \subseteq \tilde{z} + (I - |E|)^{-1} (I + |E|) \tilde{z},$$

and as in Remark 4 one gets

$$(5) \quad \|q(x^*, A^{-1}b)\| \leq \frac{2\beta_0}{1 - \beta_0} \|\text{rad } A^{-1}b\|.$$

6) Relation (5) improves Theorem 4.1 of Gay [5], who states almost the same inequality for  $\beta_0 < \frac{1}{2}$ , but with the denominator  $1 - 3\beta_0$ . However his proof seems to have a gap: In [5, p. 868, line 4] one can conclude  $w(x^i) = w(d) + w(\langle E \rangle d)$  but it is not clear whether this implies the required relation  $w(x^i) = w(d) + \sup(\|\langle E \rangle d\|)$ .

7) It would be interesting to have a result similar to Theorem 3 for preconditioned Gauss elimination.

4. On the choice of the preconditioning matrix. In Neumaier [12], the optimal choice of the preconditioning matrix  $C$  was discussed. It was shown that if  $\|I - CA\| < 1$  for some  $C \in \mathbb{R}^{n \times n}$  then  $\|I - CA\|$  takes its minimum for  $C = \tilde{A}^{-1}$ . Here we present another optimality criterion based on an upper bound for the preconditioned fixpoint inverse or Gauss inverse. We remind the reader that  $A^{-1}b$  denotes the result of interval Gauss elimination applied without pivoting to the matrix  $A$  and the right-hand side  $b$ ; cf. [12, Chap. 5].

*THEOREM 4.* Let  $A \in \mathbb{R}^{n \times n}$ ,  $b \in \mathbb{R}^n$ , and  $C \in \mathbb{R}^{n \times n}$ , and suppose that  $CA$  is an  $H$ -matrix. Then

$$(1) \quad |(CA)^F(Cb)| \subseteq (CA)^{-1}|Cb|,$$

$$(2) \quad |(CA)^G(Cb)| \subseteq (CA)^{-1}|Cb|,$$

and if  $\tilde{b} = 0$  then equality holds in (1). Moreover,  $\tilde{A}^{-1}A$  is an  $H$ -matrix and

$$(3) \quad \langle \tilde{A}^{-1}A \rangle^{-1} | \tilde{A}^{-1}b | \subseteq \langle CA \rangle^{-1} | Cb |.$$

*Proof.* By Theorem 9 of [12],  $|(CA)^F| = (CA)^{-1}$  so that  $|(CA)^F(Cb)| \cong |(CA)^F| |Cb| = (CA)^{-1} |Cb|$ , and (1) holds. If  $\beta = 0$  then  $\alpha = Cb$  satisfies  $\hat{\alpha} = 0$ , and since  $(CA)^F$  is normal [12, Chap. 4], (R8) in [12] implies  $(CA)^F(Cb) = |(CA)^F| |(Cb)| = (CA)^{-1} |Cb|$ , so that (1) holds with equality. If, instead, Theorems 3 and 4 of [12] are used one obtains in the same way relation (2).

We now observe that  $CA$  is an  $H$ -matrix, so that by [12, Thm. 5],  $\hat{A}^{-1}A$  is an  $H$ -matrix. To prove (3), let  $u > 0$  be such that  $\langle CA \rangle u > 0$ . Then  $0 \notin \langle CA \rangle_i$  for  $i = 1, \dots, n$  so that by [12, § 2.2 (O4) and § 1.3, (B9) and (B14) (the last rule contains a misprint and should read  $\rho(BA) = \rho(B)|A|$ ), we have

$$\langle CA \rangle = \langle \hat{CA} \rangle - |C| \rho(A).$$

Therefore the matrix

$$P := \langle \hat{CA} \rangle^{-1} \langle CA \rangle = I - \langle \hat{CA} \rangle^{-1} |C| \rho(A)$$

satisfies  $P \preceq I$ ,  $Pu = \langle \hat{CA} \rangle^{-1} \langle CA \rangle u > 0$ . Hence  $P$  is an  $M$ -matrix; in particular,  $P^{-1} \cong 0$ . Since

$$\begin{aligned} \langle \hat{A}^{-1}A \rangle &= \langle I + \hat{A}^{-1}\hat{A} \rangle \cong I - |\hat{A}^{-1}| \rho(A) \\ &= I - |\langle \hat{CA} \rangle^{-1} |C| \rho(A) \cong I - \langle \hat{CA} \rangle^{-1} |C| \rho(A) = P, \end{aligned}$$

we get  $P \langle \hat{A}^{-1}A \rangle^{-1} \preceq I$ ; hence  $\langle \hat{A}^{-1}A \rangle^{-1} \preceq P^{-1}$ , and finally

$$\begin{aligned} \langle \hat{A}^{-1}A \rangle^{-1} |\hat{A}^{-1}b| &\preceq P^{-1} |\hat{A}^{-1}b| = P^{-1} |\langle \hat{CA} \rangle^{-1} Cb| \\ &\preceq P^{-1} \langle \hat{CA} \rangle^{-1} |Cb| = \langle CA \rangle^{-1} |Cb|. \end{aligned} \quad \square$$

*Remarks.* 1) Inequality (3) describes a minimality property of the upper bound  $\langle CA \rangle^{-1} |Cb|$ . It would be interesting to know an optimal  $C$  for which  $|\langle CA \rangle^{-1} |Cb|$  and  $|\langle CA \rangle^G(Cb)|$  themselves take their minimum. It seems that at least for  $H$ -matrices  $A$  and  $\langle CA \rangle^G(Cb)$ , the optimal choice is not  $C = \hat{A}^{-1}$  but  $C = \hat{A}^{-1}$  where  $\hat{A} \in A$  satisfies

$$\begin{aligned} |\hat{A}_i| &= \max(|\hat{A}_i|, |\hat{A}_i|), \\ |\hat{A}_i| &= \min(|\hat{A}_i|, |\hat{A}_i|) \quad \text{for } i \neq k. \end{aligned}$$

2) A result similar to Theorem 4 is implicitly in Krawczyk [8]: If the spectral radius of  $|I - CA|$  is less than one then the expression  $(I - |I - CA|)^{-1} |Cb|$  (which is also an upper bound for  $\langle CA \rangle^G(Cb)$ ) takes its minimum for  $C = \hat{A}^{-1}$ ; indeed this follows by multiplying the third inequality in the lemma in [8] with  $(e - \hat{f})^{-1}$  which in the present notation is  $(I - |I - \hat{A}^{-1}A|)^{-1}$ .

REFERENCES

[1] G. ALFELD AND J. HERZBERGER, *Einführung in die Intervallrechnung*, Bibli. Inst., Mannheim, 1974. Expanded English translation: *Introduction to Interval Computations*, Academic Press, New York, 1983.

[2] W. BARTUŁ AND E. NUDING, *Optimale Lösung von Intervallgleichungssystemen*, *Computing*, 12 (1974), pp. 117-125.

[3] H. BELČEK, *Zur scharfen Aufgabenschätzung der Lösungsmenge bei linearen Intervallgleichungssystemen*, *Z. Angew. Math. Mech.*, 54 (1974), pp. T208-T209.

[4] P. EIJGENRAAM, *The solution of initial value problems using interval arithmetic*, *Math. Center Tracts* 144, Math. Centrum, Amsterdam, 1981.

[5] D. M. GAY, *Solving linear interval equations*, *this Journal*, 19 (1982), pp. 858-870.

[6] E. HANSEN, *Interval arithmetic in matrix comparisons*, *this Journal*, 2 (1965), pp. 308-320; 4 (1967), pp. 1-9.

[7] R. KRAWCZYK, *Newton-Algorithmen zur Bestimmung von Nullstellen mit Fehlerschranken*, *Computing*, 4 (1969), pp. 187-201.

[8] ———, *A class of interval-Newton-operators*, *Freiburger Intervall-Berichte*, 84/5 (1984), pp. 1-7.

[9] O. MAYR, *Algebraische und metrische Strukturen in der Intervallrechnung und einige Anwendungen*, *Computing*, 5 (1970), pp. 144-162.

[10] W. MEHLER, *On an interval-arithmetic matrix method*, *BIT*, 12 (1972), pp. 213-219.

[11] R. E. MOORE, *Methods and Applications of Interval Analysis*, Society for Industrial and Applied Mathematics, Philadelphia, 1979.

[12] A. NEUMAIER, *New techniques for the analysis of linear interval equations*, *Linear Algebra Appl.*, 58 (1984), pp. 273-325.

[13] ———, *Interval iteration for zeros of systems of equations*, *BIT*, 25 (1985), pp. 256-273.

[14] F. SEITZ, *Überprüfung beim Gauss-Algorithmus für lineare Gleichungssysteme*, *Freiburger Intervall-Berichte*, 84/3 (1984).

[15] P. WONGWINSIE, *Experimentelle Untersuchungen zur numerischen Auflösung von linearen Gleichungssystemen mit Fehlerfortpflanzung*, in *Interval Mathematics*, K. Nickel, ed., Springer Lecture Notes in Computer Science 29, 1975, pp. 316-325.