ON THE SYMMETRIC AND UNSYMMETRIC SOLUTION SET OF INTERVAL SYSTEMS*

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Abstract. We consider the solution set S of real linear systems Ax = b with the $n \times n$ coefficient matrix A varying between a lower bound <u>A</u> and an upper bound <u>A</u>, and with b similarly varying between <u>b</u>, <u>b</u>. First we list some properties on the shape of S if all matrices A are nonsingular. Then we restrict A to be nonsingular and symmetric deriving a complete description for the boundary of the corresponding symmetric solution set S_{sym} in the 2×2 case. Finally we derive a new criterion for the feasibility of the Cholesky method with which bounds for S_{sym} can be found.

Key words. linear interval equations, unsymmetric solution set, enclosures for the solution set of linear interval systems, symmetric linear systems, symmetric solution set, interval Cholesky method, criteria of feasibility for the interval Cholesky method

AMS subject classifications. 65F05, 65G10

1. Introduction. In [2] we introduced the interval Cholesky method in order to find an interval enclosure $[x]^C$ of the symmetric solution set

(1.1)
$$S_{\text{sym}} := \{ x \in \mathbb{R}^n | Ax = b, A = A^T \in [A], b \in [b] \},\$$

where $[A] = [A]^T$ is a given $n \times n$ matrix with real compact intervals as entries, and where [b] is a given vector with n real compact intervals as components. We showed that $[x]^C$ need not enclose the solution set

(1.2)
$$S := \{ x \in \mathbb{R}^n | Ax = b, A \in [A], b \in [b] \} \supseteq S_{\text{sym}},$$

where in this definition the symmetry of A is dropped. This phenomenon is not astonishing, since, in general, S_{sym} differs from S as was shown in [2] by a simple example.

In this paper (§4) we want to intensify our study on the symmetric solution set S_{sym} . To this end, in §3 we repeat some characteristic properties of S. Parts of them are stated and proved in [4]. We will prove them again in a much shorter way than in [4] following the lines in [8]. We then turn over to properties of S_{sym} . For 2×2 matrices S_{sym} can be represented in each orthant O as the intersection of S, O, and two sets of which the boundary is formed by conic sections. Thus, one deduces at once that in the general $n \times n$ case, the boundary ∂S_{sym} can be curvilinear in contrast to ∂S , which is shown in [4] to be the surface of a polytope.

In the second part of our paper (§5) we prove new criteria for the feasibility of the interval Cholesky method. Assuming the midpoint matrix \check{A} of [A] to be symmetric and positive definite we will show, for example, that the method results in an enclosing interval $[x]^C$ if the spectral radius of $\frac{1}{2}|\check{A}^C|d([A])$ is less than 1, where $d([A]) \in \mathbb{R}^{n \times n}$ denotes the diameter of [A] and where $|\check{A}^C|$ is a matrix which is defined later.

We mention that symmetric interval systems have also been considered by Jansson [5]. In his paper the symmetric solution set is enclosed by an iterative process.

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2. Preliminaries. We start this section with some notations that we use throughout the paper.

By \mathbb{R}^n , $\mathbb{R}^{m \times n}$, $I\mathbb{R}$, $I\mathbb{R}^n$, $I\mathbb{R}^{m \times n}$, we denote the set of real vectors with n components, the set of real $m \times n$ matrices, the set of intervals, the set of interval vectors with n components, and the set of $m \times n$ interval matrices, respectively. By "interval" we always mean a real compact interval. Interval vectors and interval matrices are vectors and matrices, respectively, with interval entries. We write intervals in brackets with the exception of degenerate intervals (so-called *point intervals*) which we identify with the element being contained, and we proceed similarly with interval vectors and interval matrix I and the null matrix O. As usual, we identify $\mathbb{R}^{n\times 1}$ and $I\mathbb{R}^{n\times 1}$ with \mathbb{R}^n and $I\mathbb{R}^n$, respectively. We use the notation $[a] = [\underline{a}, \overline{a}] \in I\mathbb{R}$ simultaneously without further reference and, in an analogous way, we write $[x] = [\underline{x}, \overline{x}] = ([x]_i) \in I\mathbb{R}^n$ and $[A] = [\underline{A}, \overline{A}] = ([a]_{ij}) \in I\mathbb{R}^{n\times n}$. For $[a], [b] \in I\mathbb{R}$ we define

	$\check{a} := (\underline{a} + \overline{a})/2$	midpoint,
	$ [a] := \max\{ \underline{a} , \overline{a} \}$	absolute value,
	$d([a]) := \overline{a} - \underline{a}$	diameter,
	$q([a],[b]) := \max\{ \underline{a} - \underline{b} , \overline{a} - \overline{b} \}$	distance,
(2.1)	$\beta([a],[b]) := [a] + q([a],[b]).$	

For interval vectors and interval matrices, these quantities are defined entrywise, i.e., they are real vectors and matrices, respectively. In particular, $|x| = (|x_i|) \in \mathbb{R}^n$ for point vectors x. We equip \mathbb{R}^n and also $\mathbb{R}^{n \times n}$ with the natural partial ordering \leq . In addition we write x < y or, equivalently, y > x for vectors $x = (x_i), y = (y_i) \in \mathbb{R}^n$ if $x_i < y_i$ for i = 1, ..., n. With the definition

$$\langle [a] \rangle := \begin{cases} 0 & \text{if } 0 \in [a] \in IR, \\ \min\{|\underline{a}|, |\overline{a}|\} & \text{otherwise,} \end{cases}$$

we construct the comparison matrix $\langle [A] \rangle := (c_{ij}) \in \mathbb{R}^{n \times n}$ of [A] by setting

$$c_{ij} := \begin{cases} \langle [a]_{ij} \rangle & \text{if } i = j, \\ -|[a]_{ij}| & \text{if } i \neq j. \end{cases}$$

We call $[A] \in IR^{n \times n}$ regular if no matrix $\tilde{A} \in [A]$ is singular, and we write $\rho(A)$ for the spectral radius of $A \in R^{n \times n}$. Intervals [a] are named zero symmetric if $\underline{a} = -\overline{a}$. For interval vectors and interval matrices zero-symmetry is defined entrywise.

We close this section by noting equivalent formulations of nonempty intersections of intervals and by recalling two properties of the function β above, which are proved in [6, Lemma 1.7.5, p. 28].

LEMMA 2.1. Let $[a], [b] \in IR$. Then the following properties are equivalent.

- (a) $[a] \cap [b] \neq \emptyset$.
- (b) $\underline{a} \leq \overline{b}$ and $\overline{a} \geq \underline{b}$.

(c) $|\check{a} - \check{b}| \leq \frac{1}{2}d([a]) + \frac{1}{2}d([b])$.

LEMMA 2.2. With β from (2.1) the following properties hold.

(a) If $[a]_i$, $[b]_i \in IR$, $[a]_i \subseteq [b]_i$ for i = 1, ..., n, then

 $\beta\left([a]_1\cdot\ldots\cdot[a]_n, \ [b]_1\cdot\ldots\cdot[b]_n\right) \leq \beta\left([a]_1, [b]_1\right)\cdot\ldots\cdot\beta\left([a]_n, [b]_n\right).$

(b) If $[a], [b] \in IR, [a] \subseteq [b] and \langle [a] \rangle > q([a], [b]), then$

$$\beta([a]^{-1}, [b]^{-1}) \le (\langle [a] \rangle - q([a], [b]))^{-1},$$

where $[c]^{-1} := \{ c^{-1} | c \in [c] \}$ for $[c] \in IR, 0 \notin [c]$.

3. The solution set S. In this section we recall some properties of the solution set S defined in (1.2). To this end, we always assume that a fixed regular interval matrix $[A] \in IR^{n \times n}$ and a fixed interval vector $[b] \in IR^n$ are given. Then the elements of S can be characterized in two equivalent ways.

THEOREM 3.1. The following three properties are equivalent.

(a) $x \in S;$

b)
$$|Ax - b| \le \frac{1}{2}d([A])|x| + \frac{1}{2}d([b])$$

(c) $[A]x \cap [b] \neq \emptyset$.

The equivalence (a) \Leftrightarrow (b) is known as Oettli–Prager criterion [7], the equivalence (a) \Leftrightarrow (c) is due to Beeck [3]. We will omit the proof.

To derive some more properties on S we decompose \mathbb{R}^n into its closed orthants $O_k, \ k = 1, \ldots, 2^n$, which are uniquely determined by the signs $s_{k_j} \in \{-1, +1\}, \ j = 1, \ldots, n$, of the components of their interior points. Hence, if O denotes some orthant, fixed by the signs s_1, \ldots, s_n , then $x = (x_i) \in O$ fulfills

(3.1)
$$x_j \begin{cases} \ge 0 & \text{if } s_j = 1, \\ \le 0 & \text{if } s_j = -1. \end{cases}$$

For [A], [b] as above, and for i, j = 1, ..., n, let

(3.2)
$$c_{ij} := \begin{cases} \underline{a}_{ij} & \text{if } s_j = 1, \\ \overline{a}_{ij} & \text{if } s_j = -1, \end{cases}$$

and

(3.3)
$$d_{ij} := \begin{cases} \overline{a}_{ij} & \text{if } s_j = 1, \\ \underline{a}_{ij} & \text{if } s_j = -1. \end{cases}$$

Denote by \underline{H}_i , \overline{H}_i , the half spaces

(3.4)
$$\begin{array}{rcl} \underline{H}_{i} & := & \left\{ y \in \mathbf{R}^{n} | \sum_{j=1}^{n} c_{ij} y_{j} \leq \overline{b}_{i} \right\} \\ \overline{H}_{i} & := & \left\{ y \in \mathbf{R}^{n} | \sum_{j=1}^{n} d_{ij} y_{j} \geq \underline{b}_{i} \right\} \end{array} \right\} \quad i = 1, \dots, n.$$

Note that \underline{H}_i , \overline{H}_i depend on the choice of the orthant O. By means of these half spaces we can represent $S \cap O$ in the following way (cf. also [8, Cor. 1.2]).

THEOREM 3.2. Let $[A] \in IR^{n \times n}$ be regular and let O denote any orthant of R^n . Then

(3.5)
$$S \cap O = \bigcap_{i=1}^{n} (\underline{H}_{i} \cap \overline{H}_{i}) \cap O.$$

In particular, if $S \cap O$ is nonempty, it is convex, compact, connected, and a polytope.

S is compact, connected, but not necessarily convex. It is the union of finitely many convex polytopes.

Proof. Let $[a] \in IR$, $\xi \in R$. Then

$$\xi \cdot [a] = \begin{cases} [\underline{\xi \underline{a}}, \underline{\xi \overline{a}}] & \text{if } \underline{\xi} \ge 0, \\ [\underline{\xi \overline{a}}, \underline{\xi \underline{a}}] & \text{if } \underline{\xi} < 0. \end{cases}$$

Hence (3.5) follows from Lemma 2.1(a), (b), from Theorem 3.1(a), (c), and from the definition of \underline{H}_i , \overline{H}_i .

Since O, \underline{H}_i , \overline{H}_i are convex, the same holds for $S \cap O$ because of (3.5). This in turn shows that $S \cap O$ is connected. The compactness and the connectivity of S follows from the same property of $[A] \times [b]$ and from the continuity of the function

$$g: \left\{ \begin{array}{rrr} [A] \times [b] & \to & \mathbf{R}^n, \\ (A,b) & \mapsto & A^{-1}b, \end{array} \right.$$

the range of which is S. Now S being compact the same holds for $S \cap O$ since O is closed. The remaining property of S follows trivially from

$$S = \bigcup_{j=1}^{2^n} (S \cap O_j)$$

and from (3.5), where O_j , $j = 1, ..., 2^n$, denote the orthants of \mathbb{R}^n numbered arbitrarily.

That S can be nonconvex is seen by the following example.

Example 3.3. Let $[A] = \begin{pmatrix} 1 & 0 \\ [-1,1] & 1 \end{pmatrix}$, $[b] = \begin{pmatrix} [-1,1] \\ 0 \end{pmatrix}$. Then S is given by $S = \{(x,y) \mid |y| \le |x| \le 1\}$ as illustrated in Fig. 1.

THEOREM 3.4. Let [A] be a point matrix. Then S is a parallelepiped; in particular, S is convex.

Proof. Let [A] = [A, A], and denote the columns of A^{-1} by c^1, \ldots, c^n . Then

$$S = \left\{ A^{-1}\underline{b} + \sum_{j=1}^{n} t_j c^j | \ 0 \le t_j \le d([b]_j), \ j = 1, \dots, n \right\}.$$

This proves the theorem.

We remark that a necessary and sufficient criterion for the convexity of S can be found in [9].

4. On the symmetric solution set S_{sym} . We now turn over to the symmetric solution set S_{sym} defined in (1.1). We again assume $[A] \in IR^{n \times n}$ to be regular, and, in addition, to fulfill

$$[A] = [A]^T,$$

which is equivalent to $\underline{A} = \underline{A}^T$ and $\overline{A} = \overline{A}^T$.

We first prove two simple properties of S_{sym} .



FIG. 1. The shape of the solution set S in Example 3.3.

THEOREM 4.1. Let $[A] = [A]^T \in IR^{n \times n}$ be regular. Then S_{sym} is compact and connected.

Proof. Define $[A]_{sym} := \{A \in [A] | A = A^T\}$. Then

(4.1)
$$f: \begin{cases} [A]_{\text{sym}} \times [b] \to \mathbb{R}^n, \\ (A,b) \mapsto A^{-1}b \end{cases}$$

is continuous. Let $\{A_k\}$ be an infinite sequence from $[A]_{\text{sym}}$. Since the (1, 1)-entries of A_k are all contained in the compact set $[a]_{11}$, there is a subsequence $\{A_k^{(1)}\}$ of $\{A_k\}$ such that its (1, 1)-entries are convergent. By the same reason one can choose a subsequence $\{A_k^{(2)}\}$ of $\{A_k^{(1)}\}$ such that the (1, 2)-entries are convergent. It is obvious that the (1, 1)-entries of $\{A_k^{(2)}\}$ keep this property. Repeating the arguments by running through the indices $(i, j), 1 \leq i \leq j \leq n$ and taking into account the symmetry of A_k shows that there is a convergent subsequence of $\{A_k\}$, which proves $[A]_{\text{sym}}$ to be compact. Therefore, $[A]_{\text{sym}} \times [b]$ is compact, and the same holds for the range S_{sym} of f.

If $A_1, A_2 \in [A]_{sym}$ then the line segment $A_1 + t(A_2 - A_1) \in [A]_{sym}, 0 \le t \le 1$. Hence $[A]_{sym}$ is connected and also $[A]_{sym} \times [b]$. Using the continuous function f from (4.1) once more shows S_{sym} to be connected. \Box

We next investigate S_{sym} in the 2 × 2 case more carefully. To this end, as in §3, we fix an orthant O given by the signs s_1, \ldots, s_n of the components of its interior

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points. We define \underline{H}_i , \overline{H}_i as in (3.2)–(3.4) and $e, f \in \mathbb{R}^n$ by

$$e_i := \begin{cases} \underline{b}_i & \text{if } s_i = 1, \\ \overline{b}_i & \text{if } s_i = -1, \end{cases}$$
$$f_i := \begin{cases} \overline{b}_i & \text{if } s_i = 1, \\ \underline{b}_i & \text{if } s_i = -1. \end{cases}$$

For n = 2 we use the sets

(4.2)
$$C^{-} := \left\{ y \in \mathbf{R}^{2} \mid \underline{a}_{11}y_{1}^{2} - \overline{a}_{22}y_{2}^{2} - f_{1}y_{1} + e_{2}y_{2} \le 0 \right\},$$

(4.3)
$$C^+ := \left\{ y \in \mathbf{R}^2 \mid \overline{a}_{11}y_1^2 - \underline{a}_{22}y_2^2 - e_1y_1 + f_2y_2 \ge 0 \right\}.$$

Obviously, each of these two sets has a conic section as boundary provided that $\underline{a}_{11}^2 + \overline{a}_{22}^2 \neq 0$ for C^- and, similarly, $\overline{a}_{11}^2 + \underline{a}_{22}^2 \neq 0$ for C^+ . As for the hyperplanes \underline{H}_i , \overline{H}_i in §3 we point out that C^- , C^+ depend on the choice of the orthant O. However, the type of the conic section is independent of O if one does not distinguish between hyperbolas and pairs of intersecting straight lines, and if one considers a single point as an ellipse. If each symmetric matrix from [A] is positive definite then $\underline{a}_{ii} > 0$, i = 1, 2, hence the boundary of C^- and C^+ is formed by hyperbolas in the above-mentioned generalized sense.

We now describe S_{sym} in the 2 × 2 case by means of S, C^- , and C^+ .

THEOREM 4.2. Let $[A] = [A]^T \in IR^{2 \times 2}$ be regular and let O denote any orthant of \mathbb{R}^2 . Then

$$(4.4) S_{\rm sym} \cap O = S \cap O \cap C^- \cap C^+$$

In particular, if $S_{sym} \cap O$ is nonempty, it is compact, but not necessarily convex.

Proof. The compactness follows from Theorem 4.1. The nonconvexity is shown by Example 4.4. It remains to prove (4.4).

 \subseteq : Let $x \in S_{sym} \cap O$. Then $x \in S \cap O$, and there exists a symmetric matrix $A \in [A]$ and a vector $b \in [b]$ such that Ax = b. With $[t] := [a]_{12} = [a]_{21}$ and $t := a_{12} = a_{21}$ we get

$$(4.5) a_{11}x_1 + tx_2 = b_1 ,$$

$$(4.6) tx_1 + a_{22}x_2 = b_2 .$$

Multiplying (4.5) by x_1 and (4.6) by x_2 and substituting tx_1x_2 we obtain

$$a_{11}x_1^2 - a_{22}x_2^2 = b_1x_1 - b_2x_2 \; .$$

Thus

(4.7)
$$x^T \begin{pmatrix} [a]_{11} & 0 \\ 0 & -[a]_{22} \end{pmatrix} x \cap x^T \begin{pmatrix} [b]_1 \\ -[b]_2 \end{pmatrix} \neq \emptyset ,$$

whence, by Lemma 2.1, we get equivalently

$$\frac{\underline{a}_{11}x_1^2 - \overline{a}_{22}x_2^2}{\overline{a}_{11}x_1^2 - \underline{a}_{22}x_2^2} \le f_1x_1 - e_2x_2 ,$$

$$\overline{a}_{11}x_1^2 - \underline{a}_{22}x_2^2 \ge e_1x_1 - f_2x_2 .$$

This means $x \in C^-$ and $x \in C^+$, respectively. Therefore, $S_{\text{sym}} \cap O \subseteq S \cap O \cap C^- \cap C^+$.

 \supseteq : Let

$$x \in S \cap O \cap C^- \cap C^+ .$$

Since $x \in S$, there are $A \in [A]$, $b \in [b]$ such that

holds. We are going to show that $A \in [A]$ in (4.9) can be chosen to be symmetric when changing $b \in [b]$ appropriately. To simplify the notation we use

$$t_1 := a_{12} \in [a]_{12}$$
 and $t_2 := a_{21} \in [a]_{21} = [a]_{12} =: [t]$

for the two off-diagonal entries of A in (4.9).

If $t_1 = t_2$ then $x \in S_{sym} \cap O$. Therefore, assume $t_1 \neq t_2$, say

(4.10)
$$t_1 < t_2$$
.

If $x_1 = 0$ then A can be replaced in (4.9) by the symmetric matrix

$$A_{\rm sym} := \left(\begin{array}{cc} a_{11} & t_1 \\ t_1 & a_{22} \end{array}\right)$$

thus showing $x \in S_{sym} \cap O$. Analogously one proceeds for $x_2 = 0$.

Let now $x_1 \neq 0$ and $x_2 \neq 0$. We first consider the case $x_1 > 0$, $x_2 > 0$, which, by (4.8), means that O is the first quadrant of \mathbb{R}^2 . Our proof is based on the equivalence of (4.9) with

(4.11)
$$t_1 = \frac{b_1 - a_{11}x_1}{x_2} \in [t], \qquad t_2 = \frac{b_2 - a_{22}x_2}{x_1} \in [t] .$$

Assume $x \notin S_{\text{sym}} \cap O$. This means that $b \in [b]$ and $A \in [A]$ from (4.9) cannot be replaced such that (4.9) is satisfied for some symmetric matrix $A_{\text{sym}} \in [A]$ and some suitably modified vector $b \in [b]$. Taking into account (4.10) we consequently obtain

(4.12)
$$\underline{t} \leq t_1 \leq t_{\max} := \frac{b_1 - \underline{a}_{11} x_1}{x_2} < t_{\min} := \frac{\underline{b}_2 - \overline{a}_{22} x_2}{x_1} \leq t_2 \leq \overline{t},$$

whence

$$\overline{b}_1 x_1 - \underline{a}_{11} x_1^2 < \underline{b}_2 x_2 - \overline{a}_{22} x_2^2 \,.$$

Since we supposed O to be the first quadrant this implies $x \notin C^-$, which contradicts (4.8).

Replacing (4.10) by $t_1 > t_2$ and assuming $x \notin S_{sym} \cap O$ yields

$$\overline{t} \ge t_1 \ge t_{\min} := \frac{\underline{b}_1 - \overline{a}_{11}x_1}{x_2} > t_{\max} := \frac{\overline{b}_2 - \underline{a}_{22}x_2}{x_1} \ge t_2 \ge \underline{t}$$

from which we get the contradiction $x \notin C^+$. Therefore,

$$(4.13) S \cap O \cap C^- \cap C^+ \subseteq S_{\text{sym}} \cap O$$

holds if O is the first quadrant O_1 .

Let now $x \in O \neq O_1$, $x_1 \neq 0$, $x_2 \neq 0$, $s_1 := sign(x_1)$, $s_2 := sign(x_2)$, $D_x := diag(s_1, s_2) \in \mathbb{R}^{2 \times 2}$. Then (4.9) is equivalent to

with $\hat{A} := D_x A D_x \in D_x[A] D_x =: [\hat{A}], \ \hat{x} := D_x x \in O_1, \ \hat{b} := D_x b \in D_x[b] =: [\hat{b}]$. Let $S, S_{\text{sym}}, C^-, C^+, e_i, f_i$ be associated with the given quantities [A], [b], and O, and let $\hat{S}, \hat{S}_{\text{sym}}, \hat{C}^-, \hat{C}^+, \hat{e}_i, \hat{f}_i$ be the corresponding quantities associated with $[\hat{A}], [\hat{b}],$ and O_1 . Since

$$s_1 f_1 = \left\{ \begin{array}{ccc} b_1 & \text{if} & s_1 = & 1 \\ -\underline{b}_1 & \text{if} & s_1 = & -1 \end{array} \right\} = \max\{s_1[b]_1\} = \hat{f}_1$$

and

$$s_2 e_2 = \left\{ \begin{array}{ccc} \underline{b}_2 & \text{if} & s_2 = & 1\\ -\overline{b}_2 & \text{if} & s_2 = & -1 \end{array} \right\} = \min\{s_2[b]_2\} = \hat{e}_2,$$

we get from $y \in C^-$ the inequality

$$0 \ge (s_1\underline{a}_{11}s_1)(s_1y_1)^2 - (s_2\overline{a}_{22}s_2)(s_2y_2)^2 - (s_1f_1)(s_1y_1) + (s_2e_2)(s_2y_2)^2 = \underline{\hat{a}}_{11}\hat{y}_1^2 - \overline{\hat{a}}_{22}\hat{y}_2^2 - f_1\hat{y}_1 + \hat{e}_2\hat{y}_2,$$

where $\hat{y} := D_x y$. Hence $y \in C^-$ implies $\hat{y} \in \hat{C}^-$, and analogously $y \in C^+$ yields $\hat{y} \in \hat{C}^+$. Therefore, $x \in S \cap O \cap C^- \cap C^+$ results in $\hat{x} \in \hat{S} \cap O_1 \cap \hat{C}^- \cap \hat{C}^+$ whence

$$(4.15) \qquad \qquad \hat{x} \in \hat{S}_{\text{sym}} \cap O_1$$

as we have proved above. Since (4.15) implies $\hat{A}_{sym}\hat{x} = \hat{b}$ for some symmetric matrix $\hat{A}_{sym} \in [\hat{A}]$ and some right-hand side $\hat{b} \in [\hat{b}]$, it yields $x \in S_{sym} \cap O$ via (4.14).

The generalization of Theorem 4.2 for the case n > 2 is not straightforward since the elimination process performed in the proof does not seem to work in this case.

Since $x \in C^- \cap C^+$ is equivalent to (4.7), we obtain immediately the subsequent corollary from Theorem 3.1(a), (c) and from Theorem 4.2.

COROLLARY 4.3. For regular matrices $[A] = [A]^T \in IR^{2\times 2}$ and $[b] \in IR^2$ the following properties are equivalent.

(a) $x \in S_{sym}$.

(b)
$$[A]x \cap [b] \neq \emptyset$$
 (i. e., $x \in S$) and
 $x^T \begin{pmatrix} [a]_{11} & 0 \\ 0 & -[a]_{22} \end{pmatrix} x \cap x^T \begin{pmatrix} [b]_1 \\ -[b]_2 \end{pmatrix} \neq \emptyset$

Note that in contrast to Theorem 4.2 no orthant enters explicitly in Corollary 4.3. Therefore, it can be viewed as an analogue of Theorem 3.1.

We now illustrate Theorem 4.2 by two examples. In particular we show that S_{sym} can be nonconvex in the orthants and that its boundary can be curvilinear.

Example 4.4. Let

$$[A] := \left(\begin{array}{cc} 5 & [-4,0] \\ [-4,0] & 5 \end{array} \right) \quad \text{and} \quad [b] := \left(\begin{array}{c} 9 \\ 0 \end{array} \right).$$

With

$$A_{t_1,t_2} := \begin{pmatrix} 5 & t_1 \\ t_2 & 5 \end{pmatrix}, \ t_1, \ t_2 \in [-4,0],$$

we get

$$A_{t_1,t_2}^{-1} = \frac{1}{25 - t_1 t_2} \begin{pmatrix} 5 & -t_1 \\ -t_2 & 5 \end{pmatrix} \ge 0$$

and

(4.16)
$$A_{t_1,t_2}^{-1} \cdot \begin{pmatrix} 9\\ 0 \end{pmatrix} = \frac{1}{25 - t_1 t_2} \begin{pmatrix} 45\\ -9t_2 \end{pmatrix}, t_1, t_2 \in [-4,0].$$

Hence S and S_{sym} are completely contained in the first quadrant O_1 . With the notations of §§3 and 4 we obtain

$$\underline{H}_1 = \{ y \in \mathbf{R}^2 \mid 5y_1 - 4y_2 \le 9 \} , \qquad \overline{H}_1 = \{ y \in \mathbf{R}^2 \mid 5y_1 \ge 9 \} , \\ \underline{H}_2 = \{ y \in \mathbf{R}^2 \mid -4y_1 + 5y_2 \le 0 \} , \qquad \overline{H}_2 = \{ y \in \mathbf{R}^2 \mid 5y_2 \ge 0 \} ,$$

hence $S = \underline{H}_1 \cap \overline{H}_1 \cap \underline{H}_2 \cap \overline{H}_2 \cap O_1$ is the triangle with the vertices (1.8, 0), (1.8, 1.44), and (5, 4). To describe S_{sym} we list the sets

$$C^{-} = \left\{ y \in \mathbf{R}^{2} \mid 5y_{1}^{2} - 5y_{2}^{2} - 9y_{1} \le 0 \right\},\$$

$$C^{+} = \left\{ y \in \mathbf{R}^{2} \mid 5y_{1}^{2} - 5y_{2}^{2} - 9y_{1} \ge 0 \right\}.$$

Then $K := C^- \cap C^+$ is the hyperbola

$$K: \quad \left(y_1 - \frac{9}{10}\right)^2 - y_2^2 = \frac{81}{100} \; .$$

By (4.16) or by Theorem 4.2 one can see that S_{sym} is that part of the right branch of K which lies between the points (1.8, 0) and (5, 4). The sets S and S_{sym} are illustrated in Fig. 2.

Our next example shows that parts of a parabola, of a circle, and straight lines can also form the boundary of S_{sym} .

Example 4.5. Let

$$[A] := \begin{pmatrix} 1 & [1,2] \\ [1,2] & [-1,0] \end{pmatrix}, \quad [b] := \begin{pmatrix} 4 \\ [1,2] \end{pmatrix}, \quad A_{\alpha,\beta,\gamma} := \begin{pmatrix} 1 & \alpha \\ \beta & \gamma \end{pmatrix} \in [A]$$

with $\alpha, \beta \in [1, 2], \gamma \in [-1, 0]$. Since det $A_{\alpha,\beta,\gamma} = \gamma - \alpha\beta \leq -1$, the interval matrix [A] is regular with

$$A_{\alpha,\beta,\gamma}^{-1} = \frac{1}{\det A_{\alpha,\beta,\gamma}} \left(\begin{array}{cc} \gamma & -\alpha \\ -\beta & 1 \end{array}\right) \ .$$

With $b_1 \geq 2b_2 \geq 2$ we get $A_{\alpha,\beta,\gamma}^{-1} \cdot b \geq 0$ for any choice $A_{\alpha,\beta,\gamma} \in [A]$, $b \in [b]$. Hence S and S_{sym} are completely contained in the first quadrant O_1 . Using the notation above we obtain for O_1 the following sets:

$$\underline{H}_1 = \{ y \in \mathbf{R}^2 \mid y_1 + y_2 \le 4 \}, \overline{H}_1 = \{ y \in \mathbf{R}^2 \mid y_1 + 2y_2 \ge 4 \}, \\ \underline{H}_2 = \{ y \in \mathbf{R}^2 \mid y_1 - y_2 \le 2 \}, \overline{H}_2 = \{ y \in \mathbf{R}^2 \mid 2y_1 \ge 1 \},$$

$$C^{-} = \left\{ y \in \mathbb{R}^{2} \mid y_{1}^{2} - 4y_{1} + y_{2} \leq 0 \right\} = \left\{ y \in \mathbb{R}^{2} \mid y_{2} \leq 4 - (y_{1} - 2)^{2} \right\},\$$

$$C^{+} = \left\{ y \in \mathbb{R}^{2} \mid y_{1}^{2} + y_{2}^{2} - 4y_{1} + 2y_{2} \geq 0 \right\} = \left\{ y \in \mathbb{R}^{2} \mid (y_{1} - 2)^{2} + (y_{2} + 1)^{2} \geq 5 \right\}.$$

The set $S = \underline{H}_1 \cap \overline{H}_1 \cap \underline{H}_2 \cap \overline{H}_2 \cap O_1$ is the convex hull of the points $(\frac{1}{2}, \frac{7}{4}), (\frac{1}{2}, \frac{7}{2}), (\frac{8}{3}, \frac{2}{3})$ and (3, 1). The boundary of $S_{\text{sym}} = S \cap O_1 \cap C^- \cap C^+$ is formed by the following four curves.



FIG. 2. The shape of the solution sets S and S_{sym} in Example 4.4.

- (i) The straight line between (¹/₂, ⁷/₄) and (⁸/₅, ⁶/₅).
 (ii) The straight line between (1, 3) and (3, 1).
 (iii) The part of the parabola y₂ = 4 (y₁ 2)² between (¹/₂, ⁷/₄) and (1, 3).
 (iv) The part of the circle (y₁ 2)² + (y₂ + 1)² = 5 between (⁸/₅, ⁶/₅) and (3, 1).

The situation is illustrated in Fig. 3.

5. Computing enclosures for S_{sym} . As was shown in [2], S_{sym} can be enclosed by the vector $[x]^C$, which results from the following interval version of the well-known Cholesky method, for which we assume $[A] = [A]^T \in IR^{n \times n}$, and $[b] \in IR^n$.

;

Step 1. " LL^T decomposition" for j := 1 to n do

$$[l]_{jj} := \left([a]_{jj} - \sum_{k=1}^{j-1} [l]_{jk}^2 \right)^{\frac{1}{2}};$$

for $i := j+1$ to n do
$$[l]_{ij} := \left([a]_{ij} - \sum_{k=1}^{j-1} [l]_{ik} [l]_{jk} \right) / [l]_{jj}$$



FIG. 3. The shape of the solution sets S and S_{sym} in Example 4.5.

Step 2. Forward substitution for i := 1 to n do

$$[y]_i := \left([b]_i - \sum_{j=1}^{i-1} [l]_{ij} [y]_j \right) / [l]_{ii} ;$$

Step 3. Backward substitution for i := n downto 1 do

$$[x]_{i}^{C} := \left([y]_{i} - \sum_{j=i+1}^{n} [l]_{ji} [x]_{j}^{C} \right) / [l]_{ii} ;$$

ICh([A], [b]) := [x]^C.

Here,

(5.1) $[a]^2 := \left\{ a^2 \mid a \in [a] \right\}$

and

$$[a]^{1/2} := \sqrt{[a]} := \{\sqrt{a} \mid a \in [a]\}$$

for intervals [a].

In contrast to the classical, i.e., noninterval Cholesky method, it is an open question when the interval Cholesky method is feasible. In [2] several criteria are given that guarantee the existence of $[x]^C$. We add here two new ones as well as a nonexistence criterion, which we formulate first. THEOREM 5.1. If $[A] = [A]^T \in IR^{n \times n}$ contains at least one symmetric matrix A which is not positive definite, then $[x]^C$ does not exist.

Proof. We first recall that a real symmetric matrix has an LL^T -decomposition with positive diagonal entries l_{ii} if and only if this matrix is positive definite (see [11]). L can be computed by the Cholesky method. Assume now that $A = A^T \in [A] = [A]^T$ is not positive definite. Then the Cholesky method will break down. This is the case if and only if for some index j either l_{jj} cannot be computed because of

$$a_{jj} - \sum_{k=1}^{j-1} l_{jk}^2 < 0$$

(see Step 1) or y_i cannot be computed because of $l_{ii} = 0$ (see Step 2). By the inclusion monotonicity of the interval arithmetic, either $[l]_{jj}$ does not exist, or $0 \in [l]_{ii}$ and the interval Cholesky method will break down.

Example 4.5 illustrates Theorem 5.1: Since

$$A = \begin{pmatrix} 1 & 1 \\ 1 & -1 \end{pmatrix} \in [A] = \begin{pmatrix} 1 & [1,2] \\ [1,2] & [-1,0] \end{pmatrix}$$

is not positive definitive, $[x]^C$ does not exist for [A]. Note, however, that the interval Gaussian algorithm is feasible for this interval matrix.

Before formulating our new feasibility criterion we need some preparations.

By Theorem 3.4 in [2] we have for [y] from Step 2 in the interval version of the Cholesky method

$$[y] = [D^n] \left([L^{n-1}] \left([D^{n-1}] \left(\dots \left([L^2] \left([D^2] \left([L^1] \left([D^1] [b] \right) \right) \right) \dots \right) \right) \right)$$

and

(5.2)
$$[x]^{C} = [D^{1}] \left([L^{1}]^{T} \left([D^{2}] \left(\dots \left([L^{n-2}]^{T} \left([D^{n-1}] \left([L^{n-1}]^{T} \left([D^{n}][y] \right) \right) \right) \dots \right) \right) \right),$$

where the diagonal matrices $[D^s]$ and the lower triangular matrices $[L^s]$ are defined for s = 1, ..., n-1 by

$$[d^{s}]_{ij} := \begin{cases} 1 & \text{if } i = j \neq s, \\ 1/[l]_{ss} & \text{if } i = j = s, \\ 0 & \text{otherwise,} \end{cases}$$
$$[l^{s}]_{ij} := \begin{cases} 1 & \text{if } i = j, \\ -[l]_{is} & \text{if } i > j = s, \\ 0 & \text{otherwise,} \end{cases}$$

with $[l]_{ij}$ from the Cholesky method. (Note that $[l]_{ij}$ is computed in the *j*th step of the " LL^T -decomposition"). By (5.2) it is easy to see that the mapping

(5.3)
$$f: \begin{cases} IR^n \to IR^n, \\ [b] \mapsto ICh([A], [b]) \end{cases}$$

is a sublinear one in the sense of [6, p. 98], i.e.,

(i)
$$[b] \subseteq [c] \Rightarrow f([b]) \subseteq f([c]),$$

(ii) $\alpha \in \mathbf{R} \Rightarrow f(\alpha[b]) = \alpha f([b]),$
(iii) $f([b] + [c]) \subseteq f([b]) + f([c]) \text{ for } [b], [c] \in \mathbf{IR}^n.$

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An easy computation yields

$$|[D^{n}]| \cdot |[L^{n-1}]| \cdot |[D^{n-1}]| \cdot \dots \cdot |[L^{2}]| \cdot |[D^{2}]| \cdot |[L^{1}]| \cdot |[D^{1}]| = \langle [L] \rangle^{-1}$$

again with $[L] = ([l]_{ij})$ from the Cholesky method. Hence, for the particular "righthand side" $[\hat{b}] := [-1, 1]e$, where $e = (1, ..., 1)^T$, one gets

$$\operatorname{ICh}([A], [\hat{b}]) = [x]^{C} = \langle [L]^{T} \rangle^{-1} \left(\langle [L] \rangle^{-1} [\hat{b}] \right) = \left(\langle [L]^{T} \rangle^{-1} \langle [L] \rangle^{-1} \right) [\hat{b}] .$$

With the abbreviation

(5.4)
$$|[A]^C| := \langle [L]^T \rangle^{-1} \langle [L] \rangle^{-1}$$

one therefore obtains for any $[b] \subseteq [\hat{b}]$ the inclusion

$$\operatorname{ICh}([A], [b]) \subseteq |[A]^C | [b] .$$

Thus, $|[A]^C|$ can be thought of as a *measure* for the width of the enclosure ICh([A], [b]) of S_{sym} that does not depend on the right-hand side [b] as long as [b] is contained in $[\hat{b}]$. The condition $[b] \subseteq [\hat{b}]$ can be considered as a sort of normalization. If it no longer holds, replace $[\hat{b}]$ by $t[\hat{b}]$ with t > 0 as small as possible such that $[b] \subseteq t[\hat{b}]$ is valid. Then

$$\operatorname{ICh}([A], [b]) \subseteq t | [A]^C | [b],$$

hence $t | [A]^C |$ is a corresponding measure.

By (5.2) we also get

$$\left| \left(\text{ ICh} \left([A], [-e^{(1)}, e^{(1)}] \right), \dots, \text{ ICh} \left([A], [-e^{(n)}, e^{(n)}] \right) \right) \right| = \left| [A]^C \right|$$

hence $|[A]^C|$ is the *absolute value* of the sublinear mapping f in the sense of [6, p. 100]. By an elementary rule of the diameter d (cf. [1]) one proves at once the property

 $d(f([b])) \ge | [A]^C | d([b])$

of f which is then called *normal* in [6, p. 102].

We next recall an equivalent definition of Step 1 in the interval Cholesky method. DEFINITION 5.2. ([2]) Let either $[A] = ([a]_{11}) \in IR^{1 \times 1}$ or

$$[A] = \begin{pmatrix} [a]_{11} & [c]^T \\ [c] & [A'] \end{pmatrix} = [A]^T \in IR^{n \times n}, \ n > 1, \ [c] \in IR^{n-1},$$
$$[A'] \in IR^{(n-1) \times (n-1)}.$$

- (a) $\Sigma_{[A]} := [A'] (1/[a]_{11}) [c][c]^T \in IR^{(n-1)\times(n-1)}$ is termed the Schur complement (of the (1,1) entry $[a]_{11}$) provided n > 1 and $0 \notin [a]_{11}$. In the product $[c][c]^T$ we assume that $[c]_i[c]_i$ is evaluated as $[c]_i^2$ (see (5.1)). $\Sigma_{[A]}$ is not defined if n = 1 or if $0 \in [a]_{11}$.
- (b) We call the pair $([L], [L]^T)$ the Cholesky decomposition of [A] if $0 < \underline{a}_{11}$ and if either n = 1 and $[L] = (\sqrt{[a]_{11}})$ or

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(5.5)
$$[L] = \begin{pmatrix} \sqrt{[a]_{11}} & 0\\ \frac{[c]}{\sqrt{[a]_{11}}} & [L'] \end{pmatrix},$$

where $([L'], ([L'])^T)$ is the Cholesky decomposition of $\Sigma_{[A]}$ provided that it exists.

As was shown in [2] the matrix [L] of the Cholesky method and that of the Definition 5.2(b) are identical.

The proof of our main result, Theorem 5.4, is heavily based on the following lemma.

LEMMA 5.3. Let the Cholesky decomposition $([L], [L]^T)$ of $[A] = [A]^T \in IR^{n \times n}$ exist, and let $[B] = [B]^T \supseteq [A]$ be such that for a suitable u > 0 we have

(5.6)
$$q([A], [B])u < \langle [L] \rangle \langle [L]^T \rangle u .$$

Then the Cholesky method is feasible for [B].

Proof by induction. The proof proceeds similarly as for Lemma 4.5.14 in [6]. Let n = 1. Then (5.6) implies u > 0. Again (5.6) together with $0 < \underline{a}_{11}$ yields

$$(\underline{a}_{11} - \underline{b}_{11})u \le q([A], [B])u < \langle [a]_{11} \rangle u = \underline{a}_{11}u ,$$

hence

 $0 < \underline{b}_{11}u$

follows. This shows $0 < \underline{b}_{11} = \langle [b]_{11} \rangle$ which proves the existence of ICh([B], [b]) for n = 1.

Assume now that the statement is true for some dimension $n \ge 1$, and let (5.6) hold for

(5.7)
$$[A] = \begin{pmatrix} [a]_{11} & [c]^T \\ [c] & [A'] \end{pmatrix} \subseteq [B] = \begin{pmatrix} [b]_{11} & [d]^T \\ [d] & [B'] \end{pmatrix} \in IR^{(n+1)\times(n+1)}$$

We first show $\underline{b}_{11} > 0$. With

(5.8)
$$q([A], [B]) = (q_{ij}) = \begin{pmatrix} q_{11} & r^T \\ r & Q' \end{pmatrix}$$

we get from (5.6)

$$\sum_{j=1}^{n+1} q_{1j} u_j < \langle [a]_{11} \rangle \, u_1 - \sum_{j=2}^{n+1} | [a]_{1j} | \, u_j \, ,$$

hence

$$\underline{a}_{11} - q_{11} = \langle [a]_{11} \rangle - q_{11} > \left\{ \sum_{j=2}^{n+1} (q_{1j} + |[a]_{1j}|) u_j \right\} / u_1 \ge 0 .$$

Together with (5.7) this implies $0 < \underline{b}_{11} = \langle [b]_{11} \rangle$, whence the Schur complement $\Sigma_{[B]} \supseteq \Sigma_{[A]}$ exists.

By our assumptions, the Schur complement $\Sigma_{[A]}$ has a Cholesky decomposition $([L'], [L']^T)$. If we can show that

(5.9)
$$q(\Sigma_{[A]}, \Sigma_{[B]})u' < \langle [L'] \rangle \langle [L']^T \rangle u'$$

holds for some vector u' > 0 then $\Sigma_{[B]}$ has a Cholesky decomposition, say $([\hat{L}'], [\hat{L}']^T)$, by the hypothesis of our induction, and with

$$[\hat{L}] := \begin{pmatrix} \sqrt{[b]_{11}} & O \\ \frac{[d]}{\sqrt{[b]_{11}}} & [\hat{L}'] \end{pmatrix},$$

we obtain the Cholesky decomposition $([\hat{L}], [\hat{L}]^T)$ of [B].

To prove (5.9) we apply β from (2.1) componentwise, and use the notation from (5.8) as well as that of Lemma 2.2. We then get

(5.10)

$$\begin{split} q(\Sigma_{[A]}, \Sigma_{[B]}) &= q\left([A'] - [c][c]^{T}[a]_{11}^{-1}, [B'] - [d][d]^{T}[b]_{11}^{-1}\right) \\ &\leq Q' + q([c][c]^{T}[a]_{11}^{-1}, [d][d]^{T}[b]_{11}^{-1}) \\ &= Q' - |[c][c]^{T}[a]_{11}^{-1}| + \beta([c][c]^{T}[a]_{11}^{-1}, [d][d]^{T}[b]_{11}^{-1}) \\ &= Q' - |[c]||[c]^{T}|\langle [a]_{11}\rangle^{-1} + \beta([c][c]^{T}[a]_{11}^{-1}, [d][d]^{T}[b]_{11}^{-1}) \\ &\leq Q' - |[c]||[c]^{T}|\langle [a]_{11}\rangle^{-1} + \beta([c], [d]) \cdot \beta([c]^{T}, [d]^{T}) \cdot \beta([a]_{11}^{-1}, [b]_{11}^{-1}) \\ &= Q' - |[c]||[c]^{T}|\langle [a]_{11}\rangle^{-1} + (|[c]| + r)(|[c]| + r)^{T}\beta([a]_{11}^{-1}, [b]_{11}^{-1}) . \end{split}$$

We now want to apply Lemma 2.2 (b) on the last factor in (5.10). To this end we must show

(5.11)
$$\langle [a]_{11} \rangle > q([a]_{11}, [b]_{11}) = q_{11}$$

Therefore, we set $u = \begin{pmatrix} u_1 \\ u' \end{pmatrix}$ in (5.6). With (5.5) and with the notation (5.8), we then obtain

$$\begin{pmatrix} q_{11} & r^T \\ r & Q' \end{pmatrix} \begin{pmatrix} u_1 \\ u' \end{pmatrix} < \begin{pmatrix} \sqrt{\langle [a]_{11} \rangle} & O \\ -\frac{|[c]|}{\sqrt{\langle [a]_{11} \rangle}} & \langle [L'] \rangle \end{pmatrix} \begin{pmatrix} \sqrt{\langle [a]_{11} \rangle} & -\frac{|[c]^T|}{\sqrt{\langle [a]_{11} \rangle}} \\ O & \langle [L']^T \rangle \end{pmatrix} \begin{pmatrix} u_1 \\ u' \end{pmatrix},$$

whence

(5.12)
$$q_{11}u_1 + r^T u' < \langle [a]_{11} \rangle u_1 - |[c]^T| u'$$

and

(5.13)
$$ru_1 + Q'u' < -|[c]|u_1 + |[c]||[c]^T|\langle [a]_{11}\rangle^{-1}u' + \langle [L']\rangle\langle [L']^T\rangle u'$$

Since $u_1 > 0$, the inequality (5.12) implies (5.11), and Lemma 2.2(b) and (5.10) yield $q(\Sigma_{[A]}, \Sigma_{[B]}) \leq Q' - |[c]| |[c]^T |\langle [a]_{11} \rangle^{-1} + (|[c]| + r)(|[c]| + r)^T (\langle [a]_{11} \rangle - q_{11})^{-1}$. Together with (5.11), (5.12), (5.13), this implies

$$q(\Sigma_{[A]}, \Sigma_{[B]})u' \leq Q'u' - |[c]| |[c]^{T}|\langle [a]_{11} \rangle^{-1}u' + (|[c]| + r)(\langle [a]_{11} \rangle - q_{11})^{-1}(|[c]| + r)^{T}u' < -ru_{1} - |[c]|u_{1} + \langle [L'] \rangle \langle [L']^{T} \rangle u' + (|[c]| + r)(\langle [a]_{11} \rangle - q_{11})^{-1}(\langle [a]_{11} \rangle u_{1} - q_{11}u_{1}) = \langle [L'] \rangle \langle [L']^{T} \rangle u'.$$

This proves (5.9) and terminates the induction.

We are now ready to prove our main result.

THEOREM 5.4. Let [A], $[B] \in IR^{n \times n}$, $[A] = [A]^T$, $[B] = [B]^T$, and suppose that ICh([A], [b]) exists. If

(5.14)
$$\rho(\mid [A]^C \mid q([A], [B])) < 1,$$

then the Cholesky method is feasible for [B].

Proof. Let Q := q([A], [B]), [C] := [A] + [-Q, Q]. Then $[B] \subseteq [C]$, and ICh([B], [b]) exists if ICh([C], [b]) does. By (5.14) the inverse of $I - |[A]^C | Q$ exists and can be represented as Neumann series

$$(I - | [A]^C | Q)^{-1} = \sum_{k=0}^{\infty} (| [A]^C | Q)^k \ge 0.$$

With any $v \in \mathbb{R}^n$ satisfying v > 0 define

(5.15)
$$u := \left(I - |[A]^C|Q\right)^{-1} |[A]^C|v.$$

Since $|[A]^C| \ge 0$ and $(I - |[A]^C|Q)^{-1} \ge 0$ are regular each of their rows contains at least one positive entry. Therefore $|[A]^C|v > 0$ and u > 0. Now (5.15) yields

$$|[A]^{C}|Qu = u - |[A]^{C}|v,$$

whence

$$Qu = \langle [L] \rangle \langle [L]^T \rangle u - v$$

$$< \langle L \rangle \langle L \rangle^T u,$$

with $([L], [L]^T)$ being the Cholesky decomposition of [A]. Hence, Lemma 5.3 guarantees the feasibility of the Cholesky method for [C] and therefore also for [B].

We illustrate Theorem 5.4 by a simple example. *Example* 5.5. Let

$$[B] := \begin{pmatrix} 4 & 2 & 2 \\ 2 & 4 & [0,2] \\ 2 & [0,2] & 4 \end{pmatrix}.$$

Then $\langle [B] \rangle \cdot (1,1,1)^T = 0$, hence $\langle [B] \rangle$ is singular. In particular, $\langle [B] \rangle$ is not an M-matrix (which requires $\langle [B] \rangle^{-1} \ge 0$; cf. [2]), whence, by definition, [B] is not an H-matrix. Therefore, Theorem 4.2 in [2] does not apply. Consider now

$$[A] := \begin{pmatrix} 4 & 2 & 2 \\ 2 & 4 & 1 \\ 2 & 1 & 4 \end{pmatrix} \subseteq [B] .$$

Since $\langle [A] \rangle$ is irreducibly diagonally dominant, the interval Cholesky method is feasible for [A] by Corollary 4.3 (ii) in [2], for example. A simple computation yields

$$[L] = \begin{pmatrix} 2 & 0 & 0 \\ 1 & \sqrt{3} & 0 \\ 1 & 0 & \sqrt{3} \end{pmatrix}, \qquad \langle [L] \rangle^{-1} = \frac{\sqrt{3}}{6} \begin{pmatrix} \sqrt{3} & 0 & 0 \\ 1 & 2 & 0 \\ 1 & 0 & 2 \end{pmatrix}$$

and

$$|[A]^{C}| = \langle [L]^{T} \rangle^{-1} \langle [L] \rangle^{-1} = \frac{1}{12} \begin{pmatrix} 5 & 2 & 2 \\ 2 & 4 & 0 \\ 2 & 0 & 4 \end{pmatrix}.$$

From

$$q([A], [B]) = \begin{pmatrix} 0 & 0 & 0 \\ 0 & 0 & 1 \\ 0 & 1 & 0 \end{pmatrix},$$

we get the matrix

$$|[A]^{C}|q([A], [B]) = \begin{pmatrix} 0 & \frac{1}{6} & \frac{1}{6} \\ 0 & 0 & \frac{1}{3} \\ 0 & \frac{1}{3} & 0 \end{pmatrix},$$

which has the eigenvalues $-\frac{1}{3}$, 0, $\frac{1}{3}$. Therefore, Theorem 5.4 applies. The elements $[\hat{l}]_{ij}$ that result from the interval Cholesky method for [B] are given by

$$[\hat{L}] = \left(\begin{array}{ccc} 2 & 0 & 0 \\ 1 & \sqrt{3} & 0 \\ 1 & [-1,1]/\sqrt{3} & [\sqrt{8},3]/\sqrt{3} \end{array} \right).$$

Our example also illustrates the following corollary.

COROLLARY 5.6. Let the midpoint matrix \check{A} of $[A] = [A]^T \in IR^{n \times n}$ be positive definite, and assume that

$$\rho\left(\frac{1}{2}|\check{A}^C| \ d([A])\right) < 1 \ .$$

Then the interval Cholesky method is feasible for [A].

Proof. Because of $[A] = [A]^T$, the matrix \check{A} is symmetric. Since it is positive definite by assumption, the interval Cholesky method is feasible for \check{A} when viewed as a point matrix. Taking into account $q(\check{A}, [A]) = \frac{1}{2}d([A])$, the assertion is a direct consequence of Theorem 5.4. \Box

6. Concluding remarks. We stress the fact that the main purpose of this paper is to give criteria for the feasibility of the interval Cholesky method. If this feasibility is guaranteed—for example, this is the case if one of the criteria presented in this paper or in [2] holds—the question arises immediately how close the symmetric solution set S_{sym} is included. Especially, what is the relation between the results of applying the Gaussian algorithm (or some other method) and the interval Cholesky method, respectively? In [2] it was shown by simple examples that generally no comparison is possible. The examples from [2] can be generalized to arbitrary large dimensions n > 2 without any difficulties. Hence up to now it is not clear under which conditions on the given interval matrix the interval Cholesky method is superior to the interval Gaussian algorithm or vice versa. The investigation of this question and/or some statistics about the width of the bounds for systems of larger dimension will be part of further research.

We also mention that for a given real system a very careful analysis of the floatingpoint Cholesky decomposition was performed in [10]. If the matrix as well as the right-hand side are afflicted with tolerances then bounds are computed for the set of all solutions for data within tolerances.

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