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# A SMOOTHING NEWTON METHOD FOR THE SECOND-ORDER CONE COMPLEMENTARITY PROBLEM 

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#### Abstract

In this paper we introduce a new smoothing function and show that it is coercive under suitable assumptions. Based on this new function, we propose a smoothing Newton method for solving the second-order cone complementarity problem (SOCCP). The proposed algorithm solves only one linear system of equations and performs only one line search at each iteration. It is shown that any accumulation point of the iteration sequence generated by the proposed algorithm is a solution to the SOCCP. Furthermore, we prove that the generated sequence is bounded if the solution set of the SOCCP is nonempty and bounded. Under the assumption of nonsingularity, we establish the local quadratic convergence of the algorithm without the strict complementarity condition. Numerical results indicate that the proposed algorithm is promising.


Keywords: second-order cone complementarity problem, smoothing function, smoothing Newton method, global convergence, quadratic convergence

MSC 2010: 90C25, 90C30, 90C33

## 1. Introduction

The second-order cone $(\mathrm{SOC})$ in $\mathbb{R}^{n}(n \geqslant 1)$, also called the Lorentz cone, is defined as

$$
\mathcal{K}^{n}:=\left\{\left(x_{1}, \bar{x}^{\mathrm{T}}\right)^{\mathrm{T}} \in \mathbb{R} \times \mathbb{R}^{n-1}: x_{1} \geqslant\|\bar{x}\|\right\}
$$

Here and below, $\|\cdot\|$ denotes the 2-norm defined by $\|x\|=\sqrt{x^{\mathrm{T}} x}$ for a vector $x$. For convenience, we write $\left(u^{\mathrm{T}}, v^{\mathrm{T}}\right)^{\mathrm{T}}$ as $(u, v)$ for any vectors $u, v \in \mathbb{R}^{n}$ throughout the paper.

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We consider the second-order cone complementarity problem (SOCCP) as follows:

$$
\begin{equation*}
\text { Find }(x, y) \in \mathbb{R}^{n} \times \mathbb{R}^{n}, \text { such that } x \in \mathcal{K}, y \in \mathcal{K}, x^{\mathrm{T}} y=0, y=F(x) \text {, } \tag{1.1}
\end{equation*}
$$

where $F: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is a continuously differentiable function, and $\mathcal{K} \subset \mathbb{R}^{n}$ is the Cartesian product of second-order cones, that is,

$$
\mathcal{K}=\mathcal{K}^{n_{1}} \times \ldots \times \mathcal{K}^{n_{r}},
$$

with $r, n_{1}, \ldots, n_{r} \geqslant 1$ and $n=\sum_{i=1}^{r} n_{i}$, and $x=\left(x^{1}, \ldots, x^{r}\right), y=\left(y^{1}, \ldots, y^{r}\right)$ with $x^{i}, y^{i} \in \mathcal{K}^{n_{i}}, i=1, \ldots, r$.

The following proposition shows that the complementarity condition on $\mathcal{K}=$ $\mathcal{K}^{n_{1}} \times \ldots \times \mathcal{K}^{n_{r}}$ can be decomposed into complementarity conditions on each $\mathcal{K}^{n_{i}}$.

Proposition 1.1 ([17]). Let $\mathcal{K}=\mathcal{K}^{n_{1}} \times \ldots \times \mathcal{K}^{n_{r}}, x=\left(x^{1}, \ldots, x^{r}\right) \in \mathbb{R}^{n_{1}} \times \ldots \times$ $\mathbb{R}^{n_{r}}$ and $y=\left(y^{1}, \ldots, y^{r}\right) \in \mathbb{R}^{n_{1}} \times \ldots \times \mathbb{R}^{n_{r}}$. Then the following relation holds:

$$
x \in \mathcal{K}, \quad y \in \mathcal{K}, x^{\mathrm{T}} y=0 \Leftrightarrow x^{i} \in \mathcal{K}^{n_{i}}, y^{i} \in \mathcal{K}^{n_{i}},\left(x^{i}\right)^{\mathrm{T}} y^{i}=0, \quad i=1, \ldots, r .
$$

In the following analysis, we assume that $\mathcal{K}=\mathcal{K}^{n}$. We do not lose any generality, because in view of Proposition 1.1, our analysis can be extended to the general case in a straightforward manner.

The SOCCP contains a wide class of problems such as the nonlinear complementarity problem (NCP) and the second-order cone programming (SOCP) [17]. A number of methods for solving the SOCCP have been proposed. They include the smoothing-regularization method [17], the derivative-free descent method [23], the merit function method (e.g., [6], [8]), the damped Gauss-Newton method [22], the nonsmooth method [7] and so on. Most of the methods are proposed for the monotone SOCCP.

Definition $1.2([17])$. The function $F: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is said to be monotone if for any $(x, y) \in \mathbb{R}^{n} \times \mathbb{R}^{n}$,

$$
(x-y)^{\mathrm{T}}(F(x)-F(y)) \geqslant 0 .
$$

Recently, smoothing Newton methods have attracted a lot of attention partially due to their encouraging convergent properties and numerical results. The main idea of this class of methods is to reformulate the problem concerned as a family of parameterized smooth equations and then to solve the smooth equations approximately by using Newton's method at each iteration. By driving the parameter to converge to zero, one can expect to find a solution to the original problem. However, in order
to obtain the local superlinear (quadratic) convergence, some algorithms (e.g., [1], [2], [5], [26]) depend on the assumptions of uniform nonsingularity and strict complementarity. Recently, Qi, Sun and Zhou [27] proposed a class of smoothing Newton methods for the NCP and box constrained variational inequalities. The Qi-Sun-Zhou method [27] was shown to be locally superlinearly/quadratically convergent without strict complementarity. Due to its simplicity and weaker assumptions imposed on smoothing functions, the Qi-Sun-Zhou method [27] has been further studied for the NCP (e.g., [14], [20], [31], [32]) and the SOCP (e.g., [10]-[12], [15], [28]).

In this paper, we present a new smoothing function which is coercive under suitable assumptions. Based on this new function, we propose a smoothing Newton algorithm for the SOCCP by modifying and extending the Qi-Sun-Zhou method [27]. It is proved that the proposed algorithm has the following nice properties:
(a) It is well-defined and a solution to the SOCCP can be obtained from any accumulation point of the iteration sequence generated by this method.
(b) If the solution set of (1.1) is nonempty and bounded, then the iteration sequence is bounded and hence it has at least one accumulation point.
(c) The algorithm solves only one system of linear equations and performs only one line search per iteration.
(d) The whole iteration sequence converges to the accumulation point globally. Furthermore, if the Jacobian of $F$ is Lipschitz continuous on $\mathbb{R}^{n}$, then the iteration sequence converges locally quadratically without strict complementarity.
The paper is organized as follows. In the next section, we introduce some preliminaries to be used in the subsequent sections. Moreover, we present the new smoothing function and give its properties. In Section 3, we propose a smoothing Newton method for solving the SOCCP and show the well-definedness of the algorithm. The global convergence and local quadratic convergence of the algorithm are investigated in Section 4. Preliminary numerical results are reported in Section 5. Some conclusions are made in Section 6.

In our notation, $\mathbb{R}^{n}$ denotes the space of $n$-dimensional real column vectors, and $\mathbb{R}_{+}^{n}\left(\right.$ or $\left.\mathbb{R}_{++}^{n}\right)$ denotes the non-negative (or positive) orthant in $\mathbb{R}^{n} . \mathbb{R}^{n_{1}} \times \ldots \times \mathbb{R}^{n_{r}}$ is identified with $\mathbb{R}^{n_{1}+\ldots+n_{r}}$. Thus, $\left(x_{1}, \ldots, x_{r}\right) \in \mathbb{R}^{n_{1}} \times \ldots \times \mathbb{R}^{n_{r}}$ is viewed as a column vector in $\mathbb{R}^{n_{1}+\ldots+n_{r}}$. $I$ denotes the identity matrix with suitable dimension. $\langle\cdot, \cdot\rangle$ denotes the Euclidean inner product. For any $x, y \in \mathbb{R}^{n}$, we write $x \succeq_{\mathcal{K}} y$ (or $x \succ_{\mathcal{K}} y$ ) if $x-y \in \mathcal{K}$ (or $x-y \in \operatorname{int} \mathcal{K}$, where int $\mathcal{K}$ denotes the interior of $\mathcal{K}$ ). For any $\alpha, \beta>0, \alpha=O(\beta)$ (or $\alpha=o(\beta)$ ) means that $\alpha / \beta$ is uniformly bounded (or tends to zero) as $\beta \rightarrow 0$.

## 2. Preliminaries and a new smoothing function

In this section, we first give a brief description of the Euclidean Jordan algebra associated with the SOC $\mathcal{K}$, which is a basic tool used in this paper. Then, we introduce a new smoothing function and give its properties.

For any vectors $x=\left(x_{1}, \bar{x}\right), s=\left(s_{1}, \bar{s}\right) \in \mathbb{R} \times \mathbb{R}^{n-1}$, their Jordan product associated with the SOC $\mathcal{K}$ is defined by

$$
x \circ s:=\left(x^{\mathrm{T}} s, x_{1} \bar{s}+s_{1} \bar{x}\right) .
$$

The identity element under this product is $\mathbf{e}:=(1,0, \ldots, 0)^{\mathrm{T}} \in \mathbb{R}^{n}$. Given an element $x=\left(x_{1}, \bar{x}\right) \in \mathbb{R} \times \mathbb{R}^{n-1}$, we define the symmetric matrix

$$
L_{x}:=\left(\begin{array}{cc}
x_{1} & \bar{x}^{\mathrm{T}} \\
\bar{x} & x_{1} I
\end{array}\right)
$$

where $I$ represents the $(n-1) \times(n-1)$ identity matrix. It is easy to verify that $x \circ s=L_{x} s=L_{s} x$ for any $x, s \in \mathbb{R}^{n}$. Moreover, $L_{x}$ is positive semidefinite (or positive definite and hence nonsingular) if and only if $x \in \mathcal{K}$ (or $x \in \operatorname{int} \mathcal{K})$.

Spectral decomposition is one of the basic concepts in Euclidean Jordan algebras. For any $x=\left(x_{1}, \bar{x}\right) \in \mathbb{R} \times \mathbb{R}^{n-1}$, its spectral decomposition with respect to the SOC $\mathcal{K}$ is defined as

$$
x=\lambda_{1} u_{1}+\lambda_{2} u_{2},
$$

in which

$$
\lambda_{i}=x_{1}+(-1)^{i}\|\bar{x}\|
$$

and

$$
u_{i}=\left\{\begin{array}{ll}
\frac{1}{2}\left(1,(-1)^{i} \frac{\bar{x}}{\|\bar{x}\|}\right) & \text { if } \bar{x} \neq 0, \\
\frac{1}{2}\left(1,(-1)^{i} \frac{\kappa}{\|\kappa\|}\right) & \text { otherwise, for any } \kappa \neq 0,
\end{array} \quad i=1,2\right.
$$

Since $x \in \mathcal{K}$ (or $x \in \operatorname{int} \mathcal{K}$ ) if and only if both $\lambda_{1}$ and $\lambda_{2}$ are nonnegative (or positive), one can define

$$
\begin{aligned}
\sqrt{x} & =\sqrt{\lambda_{1}} u_{1}+\sqrt{\lambda_{2}} u_{2} & & \text { for any } x \in \mathcal{K}, \\
x^{-1} & =\lambda_{1}^{-1} u_{1}+\lambda_{2}^{-1} u_{2} & & \text { for any } x \in \operatorname{int} \mathcal{K}, \\
x^{2} & =\lambda_{1}^{2} u_{1}+\lambda_{2}^{2} u_{2} & & \text { for any } x \in \mathbb{R}^{n} .
\end{aligned}
$$

Note that $x \circ x^{-1}=\mathbf{e}$. Moreover, for any $x \in \mathbb{R}^{n}, x^{2}=x \circ x$ and $x^{2} \in \mathcal{K}$.

For any $x \in \mathbb{R}^{n}$, we define $[x]_{+}$to be the nearest-point (in the Euclidean norm) projection of $x$ onto $\mathcal{K}$. For any $\alpha \in \mathbb{R}$, let $[\alpha]_{+}=\max \{0, \alpha\}$. Then from Proposition 3.3 in [16] we know that

$$
|x|=\sqrt{x^{2}}=\left|\lambda_{1}\right| u_{1}+\left|\lambda_{2}\right| u_{2}
$$

and the projection of $x$ onto $\mathcal{K}$ can be written as

$$
[x]_{+}=\left[\lambda_{1}\right]_{+} u_{1}+\left[\lambda_{2}\right]_{+} u_{2}=(x+|x|) / 2
$$

Define the vector-valued function by

$$
\begin{equation*}
\varphi_{0}(x, y):=2\left(x-[x-y]_{+}\right)=x+y-\sqrt{(x-y)^{2}} . \tag{2.1}
\end{equation*}
$$

In [16], it has been shown that $\varphi_{0}(x, y)$ satisfies the following property:

$$
\begin{equation*}
\varphi_{0}(x, y)=0 \Leftrightarrow x \in \mathcal{K}, \quad y \in \mathcal{K}, x^{\mathrm{T}} y=0 \tag{2.2}
\end{equation*}
$$

It is well known that $\varphi_{0}$ is typically nonsmooth, because it is not differentiable at $(0,0) \in \mathbb{R}^{n} \times \mathbb{R}^{n}$ which limits its practical applications. In order to overcome this difficulty, we can use a smoothing function of $\varphi_{0}$.

Definition 2.1 ([17]). For a nondifferentiable function $g: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$, a function $g_{\mu}: \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ with a parameter $\mu>0$ is called a smoothing function of $g$ if
(i) $g_{\mu}$ is differentiable for any $\mu>0$;
(ii) $\lim _{\mu \downarrow 0} g_{\mu}(x)=g(x)$ for any $x \in \mathbb{R}^{n}$.

In this paper, by smoothing the symmetric perturbed function of $\varphi_{0}$, we obtain a new vector-valued function $\varphi(\mu, x, y): \mathbb{R}_{+} \times \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ defined by

$$
\begin{equation*}
\varphi(\mu, x, y)=(\cos \mu+\sin \mu)(x+y)-\sqrt{(\cos \mu-\sin \mu)^{2}(x-y)^{2}+4 \mu^{2} \mathbf{e}} \tag{2.3}
\end{equation*}
$$

where $\mu$ is a real parameter. As we will show, the function $\varphi$ given in (2.3) is a smoothing function of $\varphi_{0}$ and it possesses some nice properties.

Theorem 2.2. Let $\varphi(\mu, x, y): \mathbb{R}_{+} \times \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ be defined by (2.3). Denote

$$
\begin{equation*}
\omega:=\omega(\mu, x, y)=\sqrt{(\cos \mu-\sin \mu)^{2}(x-y)^{2}+4 \mu^{2} \mathbf{e}} . \tag{2.4}
\end{equation*}
$$

Then the following results hold:
(i) $\varphi(\mu, x, y)$ is continuously differentiable at any $(\mu, x, y) \in \mathbb{R}_{++} \times \mathbb{R}^{n} \times \mathbb{R}^{n}$ with its Jacobian

$$
\varphi^{\prime}(\mu, x, y)=\left(\begin{array}{c}
(\cos \mu-\sin \mu)(x+y)-L_{\omega}^{-1}\left[-\left(1-2 \sin ^{2} \mu\right)(x-y)^{2}+4 \mu \mathbf{e}\right]  \tag{2.5}\\
(\cos \mu+\sin \mu) I-(\cos \mu-\sin \mu)^{2} L_{\omega}^{-1} L_{x-y} \\
(\cos \mu+\sin \mu) I+(\cos \mu-\sin \mu)^{2} L_{\omega}^{-1} L_{x-y}
\end{array}\right) .
$$

(ii) $\varphi(\mu, x, y)$ is a smoothing function of $\varphi_{0}(x, y)$.
$\operatorname{Proof}$. It is easy to show that $\varphi(\mu, x, y)$ is continuously differentiable at any $(\mu, x, y) \in \mathbb{R}_{++} \times \mathbb{R}^{n} \times \mathbb{R}^{n}$. For any $(\mu, x, y) \in \mathbb{R}_{++} \times \mathbb{R}^{n} \times \mathbb{R}^{n}$, it follows from (2.4) that $\omega \in \operatorname{int} \mathcal{K}$ and therefore $L_{\omega}$ is invertible. From the definition of $\omega$, we get

$$
\omega^{2}=(\cos \mu-\sin \mu)^{2}(x-y)^{2}+4 \mu^{2} \mathbf{e}
$$

Then,

$$
\begin{aligned}
& \omega_{\mu}^{\prime}=L_{\omega}^{-1}\left[-\left(1-2 \sin ^{2} \mu\right)(x-y)^{2}+4 \mu \mathbf{e}\right], \\
& \omega_{x}^{\prime}=(\cos \mu-\sin \mu)^{2} L_{\omega}^{-1} L_{x-y}, \\
& \omega_{y}^{\prime}=-(\cos \mu-\sin \mu)^{2} L_{\omega}^{-1} L_{x-y} .
\end{aligned}
$$

Therefore, due to the definition of $\varphi$, we have

$$
\begin{align*}
& \varphi_{\mu}^{\prime}(\mu, x, y)=(\cos \mu-\sin \mu)(x+y)-L_{\omega}^{-1}\left[-\left(1-2 \sin ^{2} \mu\right)(x-y)^{2}+4 \mu \mathbf{e}\right],  \tag{2.6}\\
& \varphi_{x}^{\prime}(\mu, x, y)=(\cos \mu+\sin \mu) I-(\cos \mu-\sin \mu)^{2} L_{\omega}^{-1} L_{x-y},  \tag{2.7}\\
& \varphi_{y}^{\prime}(\mu, x, y)=(\cos \mu+\sin \mu) I+(\cos \mu-\sin \mu)^{2} L_{\omega}^{-1} L_{x-y} . \tag{2.8}
\end{align*}
$$

From (2.6)-(2.8), we have the desired Jacobian formula. Now we prove (ii). For any $x=\left(x_{1}, \bar{x}\right)$ and $y=\left(y_{1}, \bar{y}\right) \in \mathbb{R} \times \mathbb{R}^{n-1}$, by the spectral factorization of $\omega^{2}$ we have

$$
\varphi(\mu, x, y)=(\cos \mu+\sin \mu)(x+y)-\left(\sqrt{\lambda_{1}(\mu)} u_{1}(\mu)+\sqrt{\lambda_{2}(\mu)} u_{2}(\mu)\right)
$$

where

$$
\begin{aligned}
& \lambda_{i}(\mu)=(\cos \mu-\sin \mu)^{2}\|x-y\|^{2}+4 \mu^{2}+2(-1)^{i}\|v(\mu)\|, \\
& u_{i}(\mu)=\left\{\begin{array}{ll}
\frac{1}{2}\left(1,(-1)^{i} \frac{v(\mu)}{\|v(\mu)\|}\right) & \text { if } v(\mu) \neq 0, \\
\frac{1}{2}\left(1,(-1)^{i} \frac{\kappa}{\|\kappa\|}\right) & \text { otherwise, for any } \kappa \neq 0,
\end{array} \quad i=1,2,\right. \\
& v(\mu)=(\cos \mu-\sin \mu)^{2}\left(x_{1}-y_{1}\right)(\bar{x}-\bar{y}) .
\end{aligned}
$$

In a similar way, we have

$$
\varphi_{0}(x, y)=x+y-\left(\sqrt{\lambda_{1}} u_{1}+\sqrt{\lambda_{2}} u_{2}\right)
$$

where

$$
\begin{aligned}
\lambda_{i} & =\|x-y\|^{2}+2(-1)^{i}\|v\|, \quad i=1,2, \\
u_{i} & =\left\{\begin{array}{ll}
\frac{1}{2}\left(1,(-1)^{i} \frac{v}{\|v\|}\right) & \text { if } v \neq 0, \\
\frac{1}{2}\left(1,(-1)^{i} \frac{\kappa}{\|\kappa\|}\right)
\end{array} \quad \text { otherwise, for any } \kappa \neq 0,\right.
\end{aligned} \quad i=1,2,
$$

Without loss of generality, we choose the same $\kappa \in \mathbb{R}^{n-1}$ as in $u_{i}(\mu)$. It is obvious that $\lim _{\mu \downarrow 0} v(\mu)=v$. Thus, we have

$$
\lim _{\mu \downarrow 0} \lambda_{i}(\mu)=\lambda_{i}, \quad \lim _{\mu \downarrow 0} u_{i}(\mu)=u_{i}, \quad i=1,2,
$$

which implies that

$$
\lim _{\mu \downarrow 0} \varphi(\mu, x, y)=\varphi_{0}(x, y)
$$

Therefore, it follows from (i) and Definition 2.1 that $\varphi(\mu, x, y)$ is a smoothing function of $\varphi_{0}(x, y)$. The proof is completed.

At the end of this section, we discuss the coerciveness of the function $\varphi$ defined by (2.3).

Theorem 2.3. Let $\varphi(\mu, x, y)$ be defined by (2.3), and $\xi, \zeta \in \mathbb{R}_{++}$with $\xi<\zeta$. Suppose that $\left\{\left(\mu_{k}, x^{k}, y^{k}\right)\right\} \subset \mathbb{R}_{++} \times \mathbb{R}^{n} \times \mathbb{R}^{n}$ is a sequence satisfying
(a) $\mu_{k} \in[\xi, \zeta],\left\{\left(x^{k}, y^{k}\right)\right\}$ is unbounded; and
(b) there is a bounded sequence $\left\{\left(u^{k}, v^{k}\right)\right\}$ such that $\left\{\left\langle x^{k}-u^{k}, y^{k}-v^{k}\right\rangle\right\}$ is bounded below.
Then $\left\{\varphi\left(\mu_{k}, x^{k}, y^{k}\right)\right\}$ is unbounded.
Proof. By using Lemma 4.2 in [19] and the fact that

$$
\begin{aligned}
\varphi= & (\cos \mu x+\sin \mu y)+(\sin \mu x+\cos \mu y) \\
& -\sqrt{[(\cos \mu x+\sin \mu y)-(\sin \mu x+\cos \mu y)]^{2}+4 \mu^{2} \mathbf{e}}
\end{aligned}
$$

we can prove the theorem similarly as Theorem 4.1 in [19]. For brevity, we omit the details here.

## 3. Description of the algorithm

The aim of this section is to propose a smoothing Newton method for solving the SOCCP. Under suitable assumptions, we show the well-definedness of our algorithm.

Let $z:=(\mu, x, y)$. We define the function $H(z): \mathbb{R}_{++} \times \mathbb{R}^{n} \times \mathbb{R}^{n} \rightarrow \mathbb{R}_{++} \times \mathbb{R}^{n} \times \mathbb{R}^{n}$ by

$$
H(z):=\left(\begin{array}{c}
\mu  \tag{3.1}\\
F(x)-y \\
\varphi(\mu, x, y)
\end{array}\right)
$$

where $\varphi$ is given by (2.3). Then, from Theorem 2.2 and (2.2), we know that $(x, y)$ is the solution to the SOCCP if and only if $H(z)=0$. Therefore, instead of solving the SOCCP, we may apply Newton's methods to solve the system of equations $H(z)=0$.

Lemma 3.1 ([30]). Let $a, b, u, v \in \mathbb{R}^{n}$ with $a \succ_{\mathcal{K}} 0, b \succ_{\mathcal{K}} 0, a \circ b \succ_{\mathcal{K}} 0$. If $\langle u, v\rangle \geqslant 0$ and $a \circ u+b \circ v=0$, then $u=v=0$.

Theorem 3.2. Let $z:=(\mu, x, y)$ and $H(z)$ be defined by (3.1). Then the following results hold:
(i) $H(z)$ is continuously differentiable at any $z:=(\mu, x, y) \in \mathbb{R}_{++} \times \mathbb{R}^{n} \times \mathbb{R}^{n}$ with its Jacobian

$$
H^{\prime}(z)=\left(\begin{array}{ccc}
1 & 0 & 0  \tag{3.2}\\
0 & F^{\prime}(x) & -I \\
\varphi_{\mu}^{\prime}(z) & \varphi_{x}^{\prime}(z) & \varphi_{y}^{\prime}(z)
\end{array}\right)
$$

where

$$
\begin{align*}
\varphi_{\mu}^{\prime}(z) & =(\cos \mu-\sin \mu)(x+y)-L_{\omega}^{-1}\left[-\left(1-2 \sin ^{2} \mu\right)(x-y)^{2}+4 \mu \mathbf{e}\right]  \tag{3.3}\\
\varphi_{x}^{\prime}(z) & =(\cos \mu+\sin \mu) I-(\cos \mu-\sin \mu)^{2} L_{\omega}^{-1} L_{x-y}  \tag{3.4}\\
\varphi_{y}^{\prime}(z) & =(\cos \mu+\sin \mu) I+(\cos \mu-\sin \mu)^{2} L_{\omega}^{-1} L_{x-y} \tag{3.5}
\end{align*}
$$

and $\omega$ is defined by (2.4).
(ii) If $F$ is a continuously differentiable and monotone function, then $H^{\prime}(z)$ is invertible for any $\mu \in(0, \pi / 2)$.

Proof. By Theorem 2.2, it is easy to show that (i) holds. Now we prove (ii). Fix any $\mu>0$. Let $\tilde{z}:=(\tilde{\mu}, \tilde{x}, \tilde{y}) \in \mathbb{R} \times \mathbb{R}^{n} \times \mathbb{R}^{n}$ be an arbitrary vector which satisfies
$H^{\prime}(z) \tilde{z}=0$. It is sufficient to prove that $\tilde{\mu}=0, \tilde{x}=0$, and $\tilde{y}=0$. From (3.2), $H^{\prime}(z) \tilde{z}=0$ gives

$$
\begin{align*}
\tilde{\mu} & =0,  \tag{3.6}\\
F^{\prime}(x) \tilde{x}-\tilde{y} & =0 \\
\varphi_{x}^{\prime}(z) \tilde{x}+\varphi_{y}^{\prime}(z) \tilde{y} & =0 .
\end{align*}
$$

By the second equation in (3.6) and the monotonicity of $F$, we have

$$
\begin{equation*}
\langle\tilde{x}, \tilde{y}\rangle=\left\langle\tilde{x}, F^{\prime}(x) \tilde{x}\right\rangle \geqslant 0 . \tag{3.7}
\end{equation*}
$$

By (3.4) and (3.5), the third equation in (3.6) yields

$$
\begin{aligned}
{[(\cos \mu} & \left.+\sin \mu) I-(\cos \mu-\sin \mu)^{2} L_{\omega}^{-1} L_{x-y}\right] \tilde{x} \\
& +\left[(\cos \mu+\sin \mu) I+(\cos \mu-\sin \mu)^{2} L_{\omega}^{-1} L_{x-y}\right] \tilde{y}=0 .
\end{aligned}
$$

Since $L_{x} y=x \circ y$ for any $x, y \in \mathbb{R}^{n}$, the above equality is equivalent to

$$
\begin{equation*}
(\cos \mu+\sin \mu)(\tilde{x}+\tilde{y})-(\cos \mu-\sin \mu)^{2} L_{\omega}^{-1}[(x-y) \circ(\tilde{x}-\tilde{y})]=0 . \tag{3.8}
\end{equation*}
$$

Premultiplying (3.8) by $L_{\omega}$ yields

$$
\begin{equation*}
(\cos \mu+\sin \mu) \omega \circ(\tilde{x}+\tilde{y})-(\cos \mu-\sin \mu)^{2}[(x-y) \circ(\tilde{x}-\tilde{y})]=0 . \tag{3.9}
\end{equation*}
$$

Since

$$
\begin{aligned}
& (\cos \mu+\sin \mu)(\tilde{x}+\tilde{y})=(\cos \mu \tilde{x}+\sin \mu \tilde{y})+(\sin \mu \tilde{x}+\cos \mu \tilde{y}), \\
& (\cos \mu-\sin \mu)(\tilde{x}-\tilde{y})=(\cos \mu \tilde{x}+\sin \mu \tilde{y})-(\sin \mu \tilde{x}+\cos \mu \tilde{y}),
\end{aligned}
$$

it follows from (3.9) that

$$
\begin{align*}
& {[\omega-(\cos \mu-\sin \mu)(x-y)] \circ[\cos \mu \tilde{x}+\sin \mu \tilde{y}]}  \tag{3.10}\\
& \quad+[\omega+(\cos \mu-\sin \mu)(x-y)] \circ[\sin \mu \tilde{x}+\cos \mu \tilde{y}]=0 .
\end{align*}
$$

Due to the definition of $\omega$, we have

$$
\omega^{2}=(\cos \mu-\sin \mu)^{2}(x-y)^{2}+4 \mu^{2} \mathbf{e} \succ_{\mathcal{K}}(\cos \mu-\sin \mu)^{2}(x-y)^{2} .
$$

Then, by Proposition 3.4 in [16], we get

$$
\begin{equation*}
\omega-(\cos \mu-\sin \mu)(x-y) \succ_{\mathcal{K}} 0, \quad \omega+(\cos \mu-\sin \mu)(x-y) \succ_{\mathcal{K}} 0 . \tag{3.11}
\end{equation*}
$$

Note that

$$
\begin{gather*}
{[\omega-(\cos \mu-\sin \mu)(x-y)] \circ[\omega+(\cos \mu-\sin \mu)(x-y)]=4 \mu^{2} \mathbf{e} \succ_{\mathcal{K}} 0,}  \tag{3.12}\\
\langle\cos \mu \tilde{x}+\sin \mu \tilde{y}, \sin \mu \tilde{x}+\cos \mu \tilde{y}\rangle \geqslant 0, \tag{3.13}
\end{gather*}
$$

where the inequality (3.13) follows from (3.7). Set

$$
\begin{array}{cl}
a:=\omega-(\cos \mu-\sin \mu)(x-y), & b:=\omega+(\cos \mu-\sin \mu)(x-y), \\
u:=\cos \mu \tilde{x}+\sin \mu \tilde{y}, & v:=\sin \mu \tilde{x}+\cos \mu \tilde{y} .
\end{array}
$$

Then, from (3.10)-(3.13) and using Lemma 3.1, we obtain that

$$
\cos \mu \tilde{x}+\sin \mu \tilde{y}=0, \quad \sin \mu \tilde{x}+\cos \mu \tilde{y}=0
$$

which yields

$$
\begin{equation*}
\cos \mu \mathbf{e} \circ \tilde{x}+\sin \mu \mathbf{e} \circ \tilde{y}=0, \quad \sin \mu \mathbf{e} \circ \tilde{x}+\cos \mu \mathbf{e} \circ \tilde{y}=0 . \tag{3.14}
\end{equation*}
$$

Since $\cos \mu>0$ and $\sin \mu>0$ for any $\mu \in(0, \pi / 2)$, we obtain that

$$
\cos \mu \mathbf{e} \succ_{\mathcal{K}} 0, \quad \sin \mu \mathbf{e} \succ_{\mathcal{K}} 0, \quad \cos \mu \mathbf{e} \circ \sin \mu \mathbf{e}=\frac{\sin 2 \mu}{2} \mathbf{e} \succ_{\mathcal{K}} 0 .
$$

Then, from (3.7) and (3.14), also using Lemma 3.1, we have $\tilde{x}=\tilde{y}=0$. Thus, the null space of $H^{\prime}(z)$ comprises only the origin, and hence $H^{\prime}(z)$ is invertible.

For any $z:=(\mu, x, y) \in \mathbb{R}_{++} \times \mathbb{R}^{n} \times \mathbb{R}^{n}$, we define the norm-function as follows:

$$
\begin{equation*}
\Psi(z):=\|H(z)\|^{2} . \tag{3.15}
\end{equation*}
$$

Now we give our smoothing Newton method.
Algorithm 3.1 (A smoothing Newton method for the SOCCP).
Step 0: Choose an accuracy parameter $\varepsilon>0$. Choose constants $\delta, \sigma \in(0,1)$, and $\mu_{0} \in(0, \pi / 2)$. Let $\left(x^{0}, y^{0}\right) \in \mathbb{R}^{n} \times \mathbb{R}^{n}$ be an arbitrary initial point. Set $\bar{z}:=\left(\mu_{0}, 0,0\right)$. Set $z^{0}:=\left(\mu_{0}, x^{0}, y^{0}\right)$. Take $\tau \in(0,1)$ such that $\mu_{0} \tau<1 / 2$ and $\tau\left\|H\left(z^{0}\right)\right\|<1$. Set $k:=0$.

Step 1: If $\left\|H\left(z^{k}\right)\right\| \leqslant \varepsilon$, then stop. Otherwise, compute

$$
\begin{equation*}
\beta_{k}:=\beta\left(z^{k}\right)=\tau \min \left\{1,\left\|H\left(z^{k}\right)\right\|\right\} \tag{3.16}
\end{equation*}
$$

Step 2: Compute $\Delta z^{k}:=\left(\Delta \mu_{k}, \Delta x^{k}, \Delta y^{k}\right) \in \mathbb{R} \times \mathbb{R}^{n} \times \mathbb{R}^{n}$ by

$$
\begin{equation*}
H^{\prime}\left(z^{k}\right) \Delta z^{k}=-H\left(z^{k}\right)+\beta_{k}\left\|H\left(z^{k}\right)\right\| \bar{z} \tag{3.17}
\end{equation*}
$$

where $H^{\prime}\left(z^{k}\right)$ denotes the Jacobian of the function $H$ at the point $z^{k}$.
Step 3: Choose $\alpha_{k}:=\delta^{l_{k}}$, where $l_{k}$ is the first nonnegative integer $l$ for which

$$
\begin{equation*}
\Psi\left(z^{k}+\delta^{l} \Delta z^{k}\right) \leqslant\left[1-\sigma\left(1-2 \mu_{0} \tau\right) \delta^{l}\right] \Psi\left(z^{k}\right) \tag{3.18}
\end{equation*}
$$

i.e., $l=0,1, \ldots$ are tried successively until the above inequality is satisfied for $l=l_{k}$.

Step 4: Set $z^{k+1}:=z^{k}+\alpha_{k} \Delta z^{k}$ and $k:=k+1$. Go to Step 1.
We note that Algorithm 3.1 solves only one system of linear equations and performs only one Armijo-type line search at each iteration. If $\left\|H\left(z^{k}\right)\right\|=0$, then $\left(x^{k}, y^{k}\right)$ is the solution to the SOCCP. So, the stopping criterion is reasonable. In the following, we give some remarks on Algorithm 3.1.

Remark 3.1. Suppose that $F$ is a continuously differentiable and monotone function and that $\left\{z^{k}\right\}$ is the iteration sequence generated by Algorithm 3.1. Then, the following results hold:
(a) From (3.18), it is easy to see that the sequence $\left\{\Psi\left(z^{k}\right)\right\}$ is monotonically decreasing, and hence, sequences $\left\{\left\|H\left(z^{k}\right)\right\|\right\}$ and $\left\{\beta_{k}\right\}$ are monotonically decreasing.
(b) Denote

$$
\begin{equation*}
\Omega:=\left\{z:=(\mu, x, y) \in \mathbb{R}_{++} \times \mathbb{R}^{n} \times \mathbb{R}^{n}: \mu \geqslant \mu_{0} \beta(z)\|H(z)\|\right\} \tag{3.19}
\end{equation*}
$$

where $\beta(\cdot)$ and $\mu_{0}$ are given in (3.16) and Step 0 of Algorithm 3.1, respectively. Then, $z^{k} \in \Omega$ for all $k \geqslant 0$. This can be obtained by induction. Since

$$
\beta\left(z^{0}\right)\left\|H\left(z^{0}\right)\right\|=\tau\left\|H\left(z^{0}\right)\right\| \min \left\{1,\left\|H\left(z^{0}\right)\right\|\right\} \leqslant \tau\left\|H\left(z^{0}\right)\right\|<1,
$$

we have $\mu_{0} \geqslant \mu_{0} \beta\left(z^{0}\right)\left\|H\left(z^{0}\right)\right\|$. So, $z^{0} \in \Omega$. Suppose that $z^{k} \in \Omega$, i.e., $\mu_{k} \geqslant$ $\mu_{0} \beta_{k}\left\|H\left(z^{k}\right)\right\|$. Then

$$
\begin{aligned}
\mu_{k+1}- & \mu_{0} \beta_{k+1}\left\|H\left(z^{k+1}\right)\right\| \\
& =\mu_{k}+\alpha_{k} \Delta \mu_{k}-\mu_{0} \beta_{k+1}\left\|H\left(z^{k+1}\right)\right\| \\
& =\left(1-\alpha_{k}\right) \mu_{k}+\alpha_{k} \mu_{0} \beta_{k}\left\|H\left(z^{k}\right)\right\|-\mu_{0} \beta_{k+1}\left\|H\left(z^{k+1}\right)\right\| \\
& \geqslant\left(1-\alpha_{k}\right) \mu_{0} \beta_{k}\left\|H\left(z^{k}\right)\right\|+\alpha_{k} \mu_{0} \beta_{k}\left\|H\left(z^{k}\right)\right\|-\mu_{0} \beta_{k+1}\left\|H\left(z^{k+1}\right)\right\| \\
& =\mu_{0}\left(\beta_{k}\left\|H\left(z^{k}\right)\right\|-\beta_{k+1}\left\|H\left(z^{k+1}\right)\right\|\right) \geqslant 0,
\end{aligned}
$$

where the second equality follows from the first equation in (3.17); the first inequality from the assumption that $z^{k} \in \Omega$; and the last inequality from the result (a). That is, $\mu_{k+1} \geqslant \mu_{0} \beta_{k+1}\left\|H\left(z^{k+1}\right)\right\|$. Thus, $z^{k} \in \Omega$ for all $k \geqslant 0$.
(c) $0<\mu_{k}<\pi / 2$ for all $k \geqslant 0$. First, we prove $\mu_{k}>0$ for all $k \geqslant 0$ by induction on $k$. Suppose that $\mu_{k}>0$ for some $k$, e.g., it is satisfied for $k=0$. Then, from the first equation in (3.17), we get for any $\alpha \in(0,1)$

$$
\mu_{k}+\alpha \Delta \mu_{k}=\mu_{k}+\alpha\left(-\mu_{k}+\mu_{0} \beta_{k}\left\|H\left(z^{k}\right)\right\|\right)=(1-\alpha) \mu_{k}+\alpha \mu_{0} \beta_{k}\left\|H\left(z^{k}\right)\right\|>0
$$

which implies that $\mu_{k+1}>0$. Thus, we have $\mu_{k}>0$ for all $k \geqslant 0$. On the other hand, from the result (b), we obtain that for all $k \geqslant 0$

$$
\mu_{k+1}=\left(1-\alpha_{k}\right) \mu_{k}+\alpha_{k} \mu_{0} \beta_{k}\left\|H\left(z^{k}\right)\right\| \leqslant\left(1-\alpha_{k}\right) \mu_{k}+\alpha_{k} \mu_{k}=\mu_{k} .
$$

Hence, the sequence $\left\{\mu_{k}\right\}$ is monotonically decreasing. This shows that $\mu_{k} \leqslant$ $\mu_{0}<\pi / 2$ holds for all $k \geqslant 0$. By Theorem 3.2, we know that $H^{\prime}\left(z^{k}\right)$ is nonsingular for all $k \geqslant 0$. Thus, the system of equations (3.17) is solvable, i.e., Step 2 is well-defined at the $k$ th iteration.
(d) For any $\alpha \in(0,1)$, denote

$$
\begin{equation*}
F_{k}(\alpha)=\Psi\left(z^{k}+\alpha \Delta z^{k}\right)-\Psi\left(z^{k}\right)-\alpha \Psi^{\prime}\left(z^{k}\right)^{\mathrm{T}} \Delta z^{k} . \tag{3.20}
\end{equation*}
$$

We obtain that $F_{k}(\alpha)=o(\alpha)$, since $\Psi(\cdot)$ is continuously differentiable around $z^{k}$. Then, from (3.17) and (3.20) we have

$$
\begin{aligned}
\Psi\left(z^{k}+\alpha \Delta z^{k}\right) & =\Psi\left(z^{k}\right)+\alpha \Psi^{\prime}\left(z^{k}\right)^{\mathrm{T}} \Delta z^{k}+F_{k}(\alpha) \\
& =\Psi\left(z^{k}\right)+2 \alpha H\left(z^{k}\right)^{\mathrm{T}} H^{\prime}\left(z^{k}\right) \Delta z^{k}+o(\alpha) \\
& =\Psi\left(z^{k}\right)+2 \alpha H\left(z^{k}\right)^{\mathrm{T}}\left[-H\left(z^{k}\right)+\beta_{k}\left\|H\left(z^{k}\right)\right\| \bar{z}\right]+o(\alpha) \\
& =(1-2 \alpha) \Psi\left(z^{k}\right)+2 \alpha \beta_{k}\left\|H\left(z^{k}\right)\right\| H\left(z^{k}\right)^{\mathrm{T}} \bar{z}+o(\alpha) \\
& \leqslant(1-\alpha) \Psi\left(z^{k}\right)+2 \alpha \mu_{0} \tau \Psi\left(z^{k}\right)+o(\alpha) \\
& =\left[1-\left(1-2 \mu_{0} \tau\right) \alpha\right] \Psi\left(z^{k}\right)+o(\alpha),
\end{aligned}
$$

where the inequality follows from the fact that $\|\bar{z}\|=\mu_{0}$ and $\beta_{k} \leqslant \tau$ by Step 0 and Step 1 of Algorithm 3.1. Since $\mu_{0} \tau<1 / 2$, there exists a constant $\bar{\alpha} \in(0,1)$ such that

$$
\Psi\left(z^{k}+\alpha \Delta z^{k}\right) \leqslant\left[1-\sigma\left(1-2 \mu_{0} \tau\right) \alpha\right] \Psi\left(z^{k}\right)
$$

holds for any $\alpha \in(0, \bar{\alpha})$ and $\sigma \in(0,1)$. This demonstrates that Step 3 is well-defined at the $k$ th iteration.

## 4. Convergence analysis

In this section, we show that any accumulation point $z^{*}$ of the iteration sequence $\left\{z^{k}\right\}$ is a solution to the system $H(z)=0$. If the accumulation point $z^{*}$ satisfies a nonsingularity assumption, then the iteration sequence $\left\{z^{k}\right\}$ converges to $z^{*}$ locally quadratically without strict complementarity. First, we prove the global convergence of Algorithm 3.1. For this purpose, we need the following result:

Lemma 4.1. Let $F$ be a continuously differentiable and monotone function and $\left\{z^{k}\right\}$ be the iteration sequence generated by Algorithm 3.1. If there exists $\mu^{*}>0$ such that $\mu_{k}>\mu^{*}$ holds for all $k \geqslant 0$ and $\left\|\left(x^{k}, y^{k}\right)\right\| \rightarrow \infty(k \rightarrow \infty)$, then

$$
\lim _{k \rightarrow \infty}\left\|H\left(z^{k}\right)\right\|=\infty
$$

Proof. By assuming that $\left\{\left\|H\left(z^{k}\right)\right\|\right\}$ is bounded, we will derive a contradiction. It follows from (3.1) that

$$
\begin{equation*}
\left\|H\left(z^{k}\right)\right\|^{2}=\mu_{k}^{2}+\left\|F\left(x^{k}\right)-y^{k}\right\|^{2}+\left\|\varphi\left(\mu_{k}, x^{k}, y^{k}\right)\right\|^{2} . \tag{4.1}
\end{equation*}
$$

Then, we obtain that $\left\{F\left(x^{k}\right)-y^{k}\right\}$ and $\left\{\varphi\left(\mu_{k}, x^{k}, y^{k}\right)\right\}$ are bounded. Denote

$$
\beta\left(x^{k}, y^{k}\right):=y^{k}-F\left(x^{k}\right) .
$$

It is easy to see that $\left\{\beta\left(x^{k}, y^{k}\right)\right\}$ is bounded and $y^{k}=\beta\left(x^{k}, y^{k}\right)+F\left(x^{k}\right)$. Suppose that $\left\{u^{k}\right\}$ is an arbitrary bounded sequence. Then $\left\{F\left(u^{k}\right)\right\}$ is bounded by the continuity of $F$. Let $v^{k}:=\beta\left(x^{k}, y^{k}\right)+F\left(u^{k}\right)$. Obviously, $\left\{v^{k}\right\}$ is bounded. Since $F$ is monotone, it follows from Definition 1.2 that

$$
\left\langle x^{k}-u^{k}, y^{k}-v^{k}\right\rangle=\left\langle x^{k}-u^{k}, F\left(x^{k}\right)-F\left(u^{k}\right)\right\rangle \geqslant 0 .
$$

Moreover, $\mu_{k} \in\left(\mu^{*}, \pi / 2\right)$. Thus, from Theorem 2.3 we have

$$
\lim _{\left\|\left(x^{k}, y^{k}\right)\right\| \rightarrow \infty}\left\|\varphi\left(\mu_{k}, x^{k}, y^{k}\right)\right\|=\infty
$$

which, together with (4.1), shows that $\left\{\left\|H\left(z^{k}\right)\right\|\right\}$ is unbounded. A contradiction is derived. The proof is completed.

Now we are in the position to give the global convergence of Algorithm 3.1. To this end, we need the following assumption:

Assumption 4.1. The solution set of the SOCCP defined by $S:=\{(x, y) \in$ $\left.\mathbb{R}^{n} \times \mathbb{R}^{n}: x \in \mathcal{K}, y \in \mathcal{K}, x^{\mathrm{T}} y=0, y=F(x)\right\}$ is nonempty and bounded.

Theorem 4.2 (Global convergence). Let $F$ be a continuously differentiable and monotone function. Suppose that $\left\{z^{k}\right\}$ is the iteration sequence generated by Algorithm 3.1. Then the following results hold:
(i) $\left\{\left\|H\left(z^{k}\right)\right\|\right\}$ is monotonically decreasing and converges to zero.
(ii) $\lim _{k \rightarrow \infty} \mu_{k}=0$. Moreover, if Assumption 4.1 holds, then $\left\{z^{k}\right\}$ is bounded.
(iii) Any accumulation point of $\left\{z^{k}\right\}$ is a solution to (1.1).

Proof. According to the result (a) in Remark 3.1, we know that $\left\{\left\|H\left(z^{k}\right)\right\|\right\}$ is monotonically decreasing and bounded from below by zero. Thus, there exists $H^{*} \geqslant$ 0 and $\beta^{*} \geqslant 0$ such that $\lim _{k \rightarrow \infty}\left\|H\left(z^{k}\right)\right\|=H^{*}$ and $\lim _{k \rightarrow \infty} \beta_{k}=\beta^{*}$. If $H^{*}=0$, then we obtain the desired result. Suppose that $H^{*}>0$. Then from (3.16) we have $\beta^{*}>0$. Since $z^{k} \in \Omega$ for all $k \geqslant 0$ by the result (b) in Remark 3.1, we obtain that $\mu_{k} \geqslant \mu_{0} \beta^{*} H^{*}>0$. Therefore, from Lemma 4.1, we know that $\left\{z^{k}:=\left(\mu_{k}, x^{k}, y^{k}\right)\right\}$ is bounded and hence it has at least one accumulation point $z^{*}:=\left(\mu^{*}, x^{*}, y^{*}\right)$. Without loss of generality, we assume that $\left\{z^{k}\right\}$ converges to $z^{*}$ as $k \rightarrow \infty$. Since $z^{*} \in \Omega$, we have $0<\mu_{0} \beta^{*} H^{*} \leqslant \mu^{*}<\pi / 2$ by the results (b) and (c) in Remark 3.1. It follows from Theorem 3.2 that $H^{\prime}\left(z^{*}\right)$ exists and is invertible. Hence, there exists a closed neighborhood $N\left(z^{*}\right)$ of $z^{*}$ such that for any $z \in N\left(z^{*}\right)$ we have $0<\mu<\pi / 2$ and $H^{\prime}(z)$ is invertible. Then, for any $z \in N\left(z^{*}\right)$, let $\Delta z:=(\Delta \mu, \Delta x, \Delta y) \in \mathbb{R} \times \mathbb{R}^{n} \times \mathbb{R}^{n}$ be the unique solution to the system of equations

$$
H^{\prime}(z) \Delta z=-H(z)+\beta(z)\|H(z)\| \bar{z}
$$

Denote

$$
g_{z}(\alpha):=\Psi(z+\alpha \Delta z)-\Psi(z)-\alpha \Psi^{\prime}(z)^{\mathrm{T}} \Delta z .
$$

Then, for any $z \in N\left(z^{*}\right)$, we have $\lim _{\alpha \rightarrow 0}\left|g_{z}(\alpha)\right| / \alpha=0$. Similarly to the result (d) in Remark 3.1, for any $\alpha \in(0,1)$ and $z \in N\left(z^{*}\right)$, we have

$$
\Psi(z+\alpha \Delta z) \leqslant\left[1-\left(1-2 \mu_{0} \tau\right) \alpha\right] \Psi(z)+o(\alpha) .
$$

Hence, we can find a positive number $\bar{\alpha} \in(0,1]$ such that

$$
\Psi(z+\alpha \Delta z) \leqslant\left[1-\sigma\left(1-2 \mu_{0} \tau\right) \alpha\right] \Psi(z)
$$

holds for any $\alpha \in(0, \bar{\alpha}], \sigma \in(0,1)$, and $z \in N\left(z^{*}\right)$. Therefore, for all sufficiently large $k$, there exists a nonnegative integer $\bar{l}$ such that $\delta^{\bar{l}} \in(0, \bar{\alpha}]$ and

$$
\Psi\left(z^{k}+\delta^{\bar{l}} \Delta z^{k}\right) \leqslant\left[1-\sigma\left(1-2 \mu_{0} \tau\right) \delta^{\bar{l}}\right] \Psi\left(z^{k}\right) .
$$

For all sufficiently large $k$, since $\alpha_{k}=\delta^{l_{k}} \geqslant \delta^{\bar{l}}$, it follows from Step 3 and Step 4 in Algorithm 3.1 that

$$
\Psi\left(z^{k+1}\right) \leqslant\left[1-\sigma\left(1-2 \mu_{0} \tau\right) \delta^{l_{k}}\right] \Psi\left(z^{k}\right) \leqslant\left[1-\sigma\left(1-2 \mu_{0} \tau\right) \delta^{\bar{l}}\right] \Psi\left(z^{k}\right)
$$

which implies that $\Psi\left(z^{k+1}\right) \leqslant C \Psi\left(z^{k}\right)$, where $C=1-\sigma\left(1-2 \mu_{0} \tau\right) \delta^{\bar{l}}<1$ is a constant, and thus $\Psi\left(z^{k}\right) \rightarrow 0$ as $k \rightarrow \infty$, i.e., $\left\|H\left(z^{k}\right)\right\| \rightarrow 0$ as $k \rightarrow \infty$. This contradicts the fact that the sequence $\left\{\left\|H\left(z^{k}\right)\right\|\right\}$ converges to $H^{*}>0$. Thus, $\left\{\left\|H\left(z^{k}\right)\right\|\right\}$ is monotonically decreasing and converges to zero.

Next, we prove (ii). Due to the definition of $H$ and $\lim _{k \rightarrow \infty}\left\|H\left(z^{k}\right)\right\|=0$, it is easy to see that $\lim _{k \rightarrow \infty} \mu_{k}=0$. From the result (b) in Remark 3.1, we know that $z^{k} \in \Omega$ for all $k \geqslant 0$. This gives

$$
\left\|H\left(z^{k}\right)\right\| \leqslant \frac{\mu_{k}}{\mu_{0} \beta\left(z^{k}\right)}
$$

Denote

$$
\Phi\left(z^{k}\right):=\binom{F\left(x^{k}\right)-y^{k}}{\varphi\left(z^{k}\right)}, \quad \lambda_{k}:=\max \left\{\frac{\mu_{k}}{\mu_{0} \tau}, \sqrt{\frac{\mu_{k}}{\mu_{0} \tau}}\right\} .
$$

Since $\beta\left(z^{k}\right)=\tau$ or $\beta\left(z^{k}\right)=\tau\left\|H\left(z^{k}\right)\right\|$, we have $\left\|\Phi\left(z^{k}\right)\right\| \leqslant\left\|H\left(z^{k}\right)\right\| \leqslant \lambda_{k}$ and $\lim _{k \rightarrow \infty} \lambda_{k}=0$. Hence, under Assumption 4.1, by using the famous mountain pass theorem (e.g., Theorem 3.5 in [32]), we can prove that $\left\{\left(x^{k}, y^{k}\right)\right\}$ is bounded similarly as Theorem 3.1 in [18]. Therefore, $\left\{z^{k}\right\}$ is bounded and hence it has at least one accumulation point $z^{*}:=\left(\mu^{*}, x^{*}, y^{*}\right)$.

Finally, we prove (iii). Let $z^{*}$ be an accumulation point of $\left\{z^{k}\right\}$. Then, there exists a subsequence, which we write without loss of generality as $\left\{z^{k}\right\}$, such that $\left\{z^{k}\right\}$ converges to $z^{*}$ as $k \rightarrow \infty$. It follows from the continuity of $H$ that $\lim _{k \rightarrow \infty}\left\|H\left(z^{k}\right)\right\|=$ $\left\|H\left(z^{*}\right)\right\|$. Thus, from the result (i) we know that $H\left(z^{*}\right)=0$ and $z^{*}$ is a solution to (1.1). So, we complete the proof.

Next, we discuss the local convergence of Algorithm 3.1. For this purpose, we need the concept of semismoothness which was originally introduced by Mifflin [21] for functionals. Qi and Sun [25] extended the definition of semismooth function to vector-valued functions. A locally Lipschitz function $F: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$, which has the generalized Jacobian $\partial F(z)$ as in Clarke [13], is said to be semismooth at $x \in \mathbb{R}^{n}$ if $\lim _{V \in \partial F\left(x+t h^{\prime}\right), h^{\prime} \rightarrow h, t \downarrow 0}\left\{V h^{\prime}\right\}$ exists for any $h \in \mathbb{R}^{n} ; F$ is said to be strongly semismooth at $x$ if $F$ is semismooth at $x$ and, for any $V \in \partial F(x+h), h \rightarrow 0$, it follows that

$$
\begin{equation*}
F(x+h)-F(x)-V h=O\left(\|h\|^{2}\right) \tag{4.2}
\end{equation*}
$$

Lemma 4.3. Let $\varphi(\mu, x, y)$ be defined by (2.3). Then, $\varphi(\mu, x, y)$ is strongly semismooth at all points $(\mu, x, y) \in \mathbb{R}_{++} \times \mathbb{R}^{n} \times \mathbb{R}^{n}$.

Proof. Theorem 4.2 in [9] shows that the function $g(x, \mu)=\sqrt{x^{2}+\mu^{2} \mathbf{e}}$ is strongly semismooth at $(x, \mu) \in \mathbb{R}^{n} \times \mathbb{R}$. By recalling the definition of $\varphi$ and the fact that the composition of strongly semismooth functions is strongly semismooth, we can obtain that $\varphi(\mu, x, y)$ is strongly semismooth at all points $(\mu, x, y) \in \mathbb{R}_{++} \times$ $\mathbb{R}^{n} \times \mathbb{R}^{n}$.

Since a vector-valued function is (strongly) semismooth if and only if its all component functions are (strongly) semismooth, by Lemma 4.3, we obtain the following corollary:

Corollary 4.4. Let $H(z)$ be defined by (3.1). Then, there exists a neighborhood $N\left(z^{*}\right)$ of $z^{*}$ such that $H(z)$ is semismooth at any point $z \in N\left(z^{*}\right)$. Furthermore, $H(z)$ is strongly semismooth at any point $z \in N\left(z^{*}\right)$ if $F^{\prime}(x)$ is locally Lipschitz continuous around $x^{*}$.

To obtain the local quadratic convergence of the algorithm, we need a few conditions. For related smoothing methods, a typical condition is that $z^{*}$ satisfies the strict complementarity condition and is nondegenerate (e.g., [3]-[5]). In this paper, we use the following assumption:

Assumption 4.2.
(i) All $V \in \partial H\left(z^{*}\right)$ are nonsingular.
(ii) $F^{\prime}(x)$ is locally Lipschitz continuous around $x^{*}$.

Theorem 4.5. Suppose that $F$ is a continuously differentiable and monotone function and that $z^{*}$ is an accumulation point of the iteration sequence $\left\{z^{k}\right\}$ generated by Algorithm 3.1. If Assumption 4.2 holds, then
(a) $z^{k+1}=z^{k}+\Delta z^{k}$ for all sufficiently large $k$,
(b) $\left\{z^{k}\right\}$ converges to $z^{*}$ quadratically, i.e., $\left\|z^{k+1}-z^{*}\right\|=O\left(\left\|z^{k}-z^{*}\right\|^{2}\right)$; moreover, $\mu_{k+1}=O\left(\mu_{k}^{2}\right)$.

Proof. Theorem 4.2 shows that $H\left(z^{*}\right)=0$. Since all $V \in \partial H\left(z^{*}\right)$ are nonsingular, it follows from Proposition 3.1 in [25] that $\left\|H^{\prime}\left(z^{k}\right)^{-1}\right\|=O(1)$ holds for all $z^{k}$ sufficiently close to $z^{*}$. Since $H(\cdot)$ is strongly semismooth around $z^{*}$, for all $z^{k}$ sufficiently close to $z^{*}$, we have

$$
\begin{equation*}
\left\|H\left(z^{k}\right)-H\left(z^{*}\right)-H^{\prime}\left(z^{k}\right)\left(z^{k}-z^{*}\right)\right\|=O\left(\left\|z^{k}-z^{*}\right\|^{2}\right) . \tag{4.3}
\end{equation*}
$$

Note that $H$ is locally Lipschitz continuous near $z^{*}$ since it is semismooth. Thus, $\left\|H\left(z^{k}\right)\right\|=O\left(\left\|z^{k}-z^{*}\right\|\right)$ holds for all $z^{k}$ sufficiently close to $z^{*}$. So, for all $z^{k}$ sufficiently close to $z^{*}$, it follows from (3.17) and (4.3) that

$$
\begin{align*}
\| z_{k}+ & \Delta z^{k}-z^{*} \|  \tag{4.4}\\
& =\left\|z_{k}+H^{\prime}\left(z^{k}\right)^{-1}\left[-H\left(z^{k}\right)+\beta_{k}\left\|H\left(z^{k}\right)\right\| z\right]-z^{*}\right\| \\
& \leqslant\left\|H^{\prime}\left(z^{k}\right)^{-1}\right\|\left[\left\|H\left(z^{k}\right)-H\left(z^{*}\right)-H^{\prime}\left(z^{k}\right)\left(z^{k}-z^{*}\right)\right\|+\mu_{0} \tau\left\|H\left(z^{k}\right)\right\|^{2}\right] \\
& =O\left(\left\|z^{k}-z^{*}\right\|^{2}\right)
\end{align*}
$$

where the inequality uses the fact that $\beta_{k}=\tau\left\|H\left(z^{k}\right)\right\|$ for all $z^{k}$ sufficiently close to $z^{*}$. Similarly as the proof of Theorem 3.1 in [24], we have $\left\|z^{k}-z^{*}\right\|=O\left(\left\|H\left(z^{k}\right)\right\|\right)$ for all $z^{k}$ sufficiently close to $z^{*}$. Hence, it follows from (4.4) that for all $z^{k}$ sufficiently close to $z^{*}$,

$$
\begin{equation*}
\left\|H\left(z^{k}+\Delta z^{k}\right)\right\|=O\left(\left\|z^{k}+\Delta z^{k}-z^{*}\right\|\right)=O\left(\left\|z^{k}-z^{*}\right\|^{2}\right)=O\left(\left\|H\left(z^{k}\right)\right\|^{2}\right) \tag{4.5}
\end{equation*}
$$

From Theorem 4.2, we know that $\left\|H\left(z^{k}\right)\right\| \rightarrow 0$ as $k \rightarrow \infty$. Thus, (4.5) implies that $\alpha_{k}=1$ for all $z^{k}$ sufficiently close to $z^{*}$. This, together with (4.4), indicates that the result (a) holds, and for all $z^{k}$ sufficiently close to $z^{*}$,

$$
\left\|z^{k+1}-z^{*}\right\|=O\left(\left\|z^{k}-z^{*}\right\|^{2}\right)
$$

Next, we prove the second part of (b). It follows from (a) that for all sufficiently large $k$,

$$
\begin{equation*}
\mu_{k+1}=\mu_{k}+\Delta \mu_{k}=\mu_{0} \beta_{k}\left\|H\left(z^{k}\right)\right\|=\mu_{0} \tau\left\|H\left(z^{k}\right)\right\|^{2} \tag{4.6}
\end{equation*}
$$

Thus, for all sufficiently large $k$, by using (4.5) and (4.6) we have

$$
\frac{\mu_{k+1}}{\mu_{k}^{2}}=\frac{\left\|H\left(z^{k}\right)\right\|^{2}}{\mu_{0} \tau\left\|H\left(z^{k-1}\right)\right\|^{4}}=\frac{O\left(\left\|H\left(z^{k-1}\right)\right\|^{4}\right)}{\mu_{0} \tau\left\|H\left(z^{k-1}\right)\right\|^{4}}
$$

which implies that $\mu_{k+1}=O\left(\mu_{k}^{2}\right)$. This completes the proof.

## 5. Numerical experiments

In this section, we conduct some numerical experiments to solve some SOCCP's by Algorithm 3.1 and report the computational results. All experiments were performed on a personal computer with 2.0 GB memory and $\operatorname{Intel}(\mathrm{R})$ Pentium( R$)$ Dual-Core CPU $2.93 \mathrm{GHz} \times 2$. The operating system was Windows XP and the computer codes were all written in Matlab 7.0.1. In all the experiments, we used $\left\|H\left(z^{k}\right)\right\| \leqslant 10^{-8}$ as the stopping criterion.

### 5.1. Linear case

Find $(x, y) \in \mathbb{R}^{n} \times \mathbb{R}^{n}$ such that

$$
\begin{equation*}
x \in \mathcal{K}, \quad y \in \mathcal{K}, \quad x^{\mathrm{T}} y=0, \quad y=M x+q, \tag{5.1}
\end{equation*}
$$

where $M \in \mathbb{R}^{n \times n}$ and $q \in \mathbb{R}^{n}$. We used $x^{0}=\mathbf{e}$ and $y^{0}=0$ as the starting points.
In the first set of experiments, we set

$$
M=\operatorname{diag}(1 / n, 2 / n, \ldots, 1), \quad q=(-1,-1, \ldots,-1)^{\mathrm{T}} .
$$

The parameters used in Algorithm 3.1 were chosen as $\mu_{0}=0.1, \sigma=0.5, \delta=0.8$, $\tau=0.95 /\left(1+\left\|H\left(z^{0}\right)\right\|\right)$. The results are listed in Tab. 1, where $n$ denotes the problem size; IT and CPU denote the number of iterations and the CPU time in seconds, respectively; FV and MU denote the value of $\left\|H\left(z^{k}\right)\right\|$ and $\mu_{k}$ when the algorithm terminates.

| $n$ | IT | CPU | FV | MU |
| ---: | ---: | :---: | :--- | :--- |
| 8 | 6 | 0.015 | $1.6874 \times 10^{-9}$ | $1.4683 \times 10^{-9}$ |
| 16 | 8 | 0.015 | $4.2556 \times 10^{-11}$ | $2.8190 \times 10^{-14}$ |
| 32 | 9 | 0.016 | $4.8565 \times 10^{-10}$ | $6.5322 \times 10^{-14}$ |
| 64 | 11 | 0.016 | $2.5334 \times 10^{-9}$ | $2.3882 \times 10^{-10}$ |
| 128 | 15 | 0.531 | $6.1484 \times 10^{-14}$ | $1.5337 \times 10^{-18}$ |
| 256 | 21 | 5.516 | $4.1897 \times 10^{-14}$ | $2.9229 \times 10^{-17}$ |

Table 1. Numerical results for the linear SOCCP of various problem sizes.

In the second set of experiments, the matrix $M$ was obtained by setting $M=$ $N^{\mathrm{T}} N$, where $N$ was a square matrix. Elements of $N$ and $q$ were chosen randomly from the interval $[0,1]$. The random problems of each case were generated 10 times. Tab. 2 gives the results when we chose parameters in Algorithm 3.1 as $\mu_{0}=0.1$, $\sigma=0.5, \delta=0.8, \tau=0.95 /\left(1+\left\|H\left(z^{0}\right)\right\|\right)$, in which $n$ denotes the problem size;

MIT and MCPU denote the maximum values of the number of iterations and the CPU time in seconds, respectively; AIT and ACPU denote the average values of the number of iterations and the CPU time in seconds, respectively; MFV and MGAP denote the maximum values of $\|H(z)\|$ and $\left|x^{\mathrm{T}} y\right|$ at the final iteration in the 10 trials, respectively.

| $n$ | MIT | AIT | MCPU | ACPU | MFV | MGAP |
| :---: | ---: | :---: | :---: | :---: | :--- | :--- |
| 100 | 7 | 6.4 | 0.15 | 0.12 | $8.0853 \times 10^{-9}$ | $3.9338 \times 10^{-9}$ |
| 200 | 9 | 7.3 | 1.05 | 0.85 | $2.0309 \times 10^{-10}$ | $5.2779 \times 10^{-11}$ |
| 300 | 8 | 7.8 | 3.00 | 2.92 | $2.3614 \times 10^{-9}$ | $9.5040 \times 10^{-10}$ |
| 400 | 9 | 8.5 | 7.63 | 7.16 | $7.6918 \times 10^{-9}$ | $2.6244 \times 10^{-9}$ |
| 500 | 10 | 8.8 | 16.03 | 14.11 | $4.6039 \times 10^{-9}$ | $1.1348 \times 10^{-9}$ |
| 600 | 9 | 8.6 | 24.38 | 23.28 | $7.4057 \times 10^{-10}$ | $2.2742 \times 10^{-10}$ |
| 700 | 9 | 8.8 | 38.19 | 37.26 | $7.5673 \times 10^{-10}$ | $1.9741 \times 10^{-10}$ |
| 800 | 12 | 9.4 | 73.31 | 55.65 | $7.8235 \times 10^{-9}$ | $1.5626 \times 10^{-9}$ |

Table 2. Numerical results for the linear SOCCP of various problem sizes.
Table 3 gives the results for our second set of experiments when we chose parameters in Algorithm 3.1 as $\mu_{0}=0.1, \sigma=0.25, \delta=0.5, \tau=0.5 /\left(2+\left\|H\left(z^{0}\right)\right\|\right)$.

| $n$ | MIT | AIT | MCPU | ACPU |
| :---: | :---: | :---: | :---: | :---: |
| 100 | 7 | 6.4 | 0.14 | 0.12 |
| 200 | 8 | 7.2 | 0.94 | 0.84 |
| 300 | 8 | 7.3 | 3.00 | 2.68 |
| 400 | 9 | 7.9 | 7.65 | 6.71 |
| 500 | 9 | 8.2 | 14.59 | 13.26 |
| 600 | 9 | 8.1 | 24.56 | 22.11 |
| 700 | 10 | 8.5 | 42.77 | 36.34 |
| 800 | 12 | 9.2 | 75.50 | 54.89 |

Table 3. Numerical results for the linear SOCCP of various problem sizes.
From the results in Tabs. 2 and 3, we may see that Algorithm 3.1 is effective for solving the SOCCP (5.1). It can solve all the test problems and can deal with largescale SOCCP problems. It can find a solution point meeting the desired accuracy in few iterations and in short CPU time. In addition, we also find that there are slight changes in results for different values of $\sigma, \delta$, and $\tau$.

In our third set of experiments, elements of $q$ were chosen randomly from the interval $[-1,1]$ and $M=N^{\mathrm{T}} N$, where $N \in \mathbb{R}^{n \times n}$ was a sparse matrix (the density of $N$ is nonzero) whose elements were chosen randomly from the interval $[0,1]$. The
random problems were generated 10 times for each nonzero density with problem size $n=100$ and $n=500$. The parameters used in Algorithm 3.1 were chosen as $\mu_{0}=0.1, \sigma=0.5, \delta=0.8, \tau=0.95 /\left(1+\left\|H\left(z^{0}\right)\right\|\right)$. The tested results are listed in Tabs. 4 and 5, where Dens. denotes the nonzero density of matrix $N$; AIT and ACPU denote the average values of the number of iterations and the CPU time in seconds, respectively; MFV and MMU denote the maximum values of $\|H(z)\|$ and $\mu$ at the final iteration in the 10 trials for each nonzero density, respectively.

| Dens. (\%) | AIT | ACPU | MFV | MMU |
| :---: | :---: | :---: | :---: | :---: |
| 5 | 7.6 | 0.217 | $6.0586 \times 10^{-10}$ | $3.5705 \times 10^{-12}$ |
| 10 | 7.8 | 0.245 | $1.5797 \times 10^{-9}$ | $2.2005 \times 10^{-11}$ |
| 20 | 7.1 | 0.224 | $5.2156 \times 10^{-9}$ | $7.3635 \times 10^{-11}$ |
| 40 | 7.2 | 0.241 | $7.6694 \times 10^{-9}$ | $4.6839 \times 10^{-11}$ |
| 60 | 7.2 | 0.209 | $5.7693 \times 10^{-9}$ | $1.9911 \times 10^{-11}$ |
| 80 | 7.3 | 0.236 | $3.0560 \times 10^{-9}$ | $3.7419 \times 10^{-12}$ |

Table 4. Numerical results for the linear SOCCP with different nonzero density ( $n=100$ ).

| Dens. (\%) | AIT | ACPU | MFV | MMU |
| :---: | :---: | :---: | :---: | :---: |
| 5 | 7.9 | 23.21 | $1.0290 \times 10^{-9}$ | $2.6611 \times 10^{-11}$ |
| 15 | 8.0 | 23.87 | $3.9807 \times 10^{-9}$ | $7.5880 \times 10^{-12}$ |
| 35 | 8.2 | 23.18 | $8.5395 \times 10^{-9}$ | $4.1456 \times 10^{-12}$ |
| 55 | 8.2 | 22.37 | $9.2832 \times 10^{-9}$ | $1.5064 \times 10^{-12}$ |
| 75 | 7.8 | 20.47 | $8.7656 \times 10^{-10}$ | $1.3906 \times 10^{-13}$ |
| 95 | 8.3 | 22.49 | $2.4070 \times 10^{-9}$ | $1.9853 \times 10^{-13}$ |

Table 5. Numerical results for the linear SOCCP with different nonzero density ( $n=500$ ).

The results in Tabs. 4 and 5 show that our algorithm performs well. It has good convergence and numerical stability. The number of iterations and CPU time slightly change with the sparsity of $N$. We have also tested some other problems, and the computation effect is similar.

### 5.2. Nonlinear case

Find $(x, y) \in \mathbb{R}^{5} \times \mathbb{R}^{5}$ such that

$$
\begin{equation*}
x \in \mathcal{K}, \quad y \in \mathcal{K}, \quad x^{\mathrm{T}} y=0, \quad y=F(x) \tag{5.2}
\end{equation*}
$$

where $\mathcal{K}=\mathcal{K}^{3} \times \mathcal{K}^{2}$ and $F: \mathbb{R}^{5} \rightarrow \mathbb{R}^{5}$ is given by

$$
F(x)=\left(\begin{array}{c}
24\left(2 x_{1}-x_{2}\right)^{3}+\exp \left(x_{1}-x_{3}\right)-4 x_{4}+x_{5} \\
-12\left(2 x_{1}-x_{2}\right)^{3}+3\left(3 x_{2}+5 x_{3}\right) / \sqrt{1+\left(3 x_{2}+5 x_{3}\right)^{2}}-6 x_{4}-7 x_{5} \\
-\exp \left(x_{1}-x_{3}\right)+5\left(3 x_{2}+5 x_{3}\right) / \sqrt{1+\left(3 x_{2}+5 x_{3}\right)^{2}}-3 x_{4}+5 x_{5} \\
4 x_{1}+6 x_{2}+3 x_{3}-1 \\
-x_{1}+7 x_{2}-5 x_{3}+2
\end{array}\right) .
$$

From the analysis in [17], we know that the function $F(x)$ is monotone.
In the experiments, the initial points of Algorithm 3.1 were chosen randomly. We set the parameters in Algorithm 3.1 as $\sigma=0.5, \delta=0.8, \tau=0.95 /\left(1+\left\|H\left(z^{0}\right)\right\|\right)$ and $\mu_{0}$ was a random number in $[1,10]$. We tested 10 times for this problem. Numerical results are summarized in Tab. 6, where IT and CPU denote the number of iterations and the CPU time in seconds, respectively; FV and GAP denote the values of $\|H(z)\|$ and $\left|x^{\mathrm{T}} y\right|$ at the final iteration, respectively.

| $\mu_{0}$ | IT | CPU | FV | GAP |
| :---: | :---: | :---: | :---: | :---: |
| 7.3927 | 12 | 0.017 | $2.1610 \times 10^{-12}$ | $4.1943 \times 10^{-13}$ |
| 8.2521 | 20 | 0.016 | $3.0839 \times 10^{-11}$ | $3.3006 \times 10^{-12}$ |
| 1.7885 | 12 | 0.015 | $1.5593 \times 10^{-15}$ | $1.3878 \times 10^{-17}$ |
| 3.7678 | 10 | 0.016 | $1.2750 \times 10^{-10}$ | $2.6257 \times 10^{-11}$ |
| 7.4862 | 11 | 0.015 | $3.1223 \times 10^{-12}$ | $5.8269 \times 10^{-13}$ |
| 4.8751 | 13 | 0.016 | $1.8743 \times 10^{-15}$ | $1.3878 \times 10^{-16}$ |
| 7.3947 | 11 | 0.017 | $4.7698 \times 10^{-9}$ | $8.9623 \times 10^{-10}$ |
| 2.7836 | 12 | 0.016 | $3.5686 \times 10^{-14}$ | $6.0091 \times 10^{-15}$ |
| 7.1534 | 12 | 0.015 | $1.6149 \times 10^{-12}$ | $2.7273 \times 10^{-13}$ |
| 6.6358 | 20 | 0.015 | $9.8931 \times 10^{-10}$ | $1.7717 \times 10^{-10}$ |

Table 6. Numerical results for the nonlinear SOCCP with various initial points.

### 5.3. Comparison with interior point method

We consider the following problem:

$$
\begin{array}{ll}
\min & \sum_{i=1}^{r} c_{i}^{\mathrm{T}} x_{i} \\
\text { s.t. } & \sum_{i=1}^{r} A_{i} x_{i}=b, \\
& x_{i} \in \mathcal{K}^{n_{i}}, \quad i=1, \ldots, r
\end{array}
$$

where $c_{i}, x_{i} \in \mathbb{R}^{n_{i}}, A_{i} \in \mathbb{R}^{m \times n_{i}}$, and $\mathcal{K}^{n_{i}}$ is an $n_{i}$-dimensional second-order cone. Let $y \in \mathbb{R}^{m}$ be the dual variable of the equality constraint. The KKT system of
the SOCP is

$$
\begin{equation*}
\sum_{i=1}^{r} A_{i} x_{i}=b, \quad x_{i}^{\mathrm{T}}\left(c_{i}-A_{i}^{\mathrm{T}} y\right)=0, \quad x_{i} \in \mathcal{K}^{n_{i}}, \quad c_{i}-A_{i}^{\mathrm{T}} y \in \mathcal{K}^{n_{i}}, \quad i=1, \ldots, r \tag{5.3}
\end{equation*}
$$

Let $n=\sum_{i=1}^{r} n_{i}, c=\left(c_{1}, \ldots, c_{r}\right), x=\left(x_{1}, \ldots, x_{r}\right), A=\left(A_{1}, \ldots, A_{r}\right)$. We solved the problem (5.3) by Algorithm 3.1 and SDPT3 [29], a successful interior point method software for the SOCP. The tested problems were randomly generated with sizes $n(=2 m)$ from 100 to 400 with each $n_{i}=5$. We generated a random matrix $A$ and a random vector $x$ in the second-order cone which give a right-hand side $b=A x$ and hence the problem is feasible. Moreover, we generated a random vector $c$ in the second-order cone, so the optimal value of the problem is obtainable. Throughout the computational experiments, the parameters used in Algorithm 3.1 were chosen as $\mu_{0}=2 \times 10^{-3}, \delta=0.65, \sigma=0.05, \tau=0.95 /\left(1+\left\|H\left(z^{0}\right)\right\|\right)$. Denote $\mathbf{e}^{n_{i}}=$ $(1,0, \ldots, 0)^{\mathrm{T}}$, an $n_{i}$-dimensional vector; and $\mathbf{e}=\left(\mathbf{e}^{n_{1}}, \ldots, \mathbf{e}^{n_{r}}\right)$. We chose $x^{0}=\mathbf{e}$, $y^{0}=0$ as the starting points.

The random problems of each case were generated 5 times. Tab. 7 and Tab. 8 give the numerical results when we implemented Algorithm 3.1 and SDPT3 for the problem (5.3), respectively, where mIT denotes the minimum value of the number of iterations, AIT denotes the average value of the number of iterations, MCPU denotes the maximum value of the CPU time in seconds, ACPU denotes the average value of the CPU time in seconds, AFV denotes the average value of $\left\|H\left(z^{k}\right)\right\|$ when the algorithm terminates among the 5 tests, and $\mathbf{m F V}$ denotes the minimum value of $\left\|H\left(z^{k}\right)\right\|$ when the algorithm terminates among the 5 tests. Comparing Tab. 7 and Tab. 8, we may see that the numerical results of Algorithm 3.1 are better than those by SDPT3, either in terms of the number of iterations or the CPU time.

| $m$ | $n$ | mIT | AIT | MCPU | ACPU | AFV | mFV |
| ---: | :---: | :---: | :---: | :---: | :---: | :--- | :---: |
| 50 | 100 | 11 | 12.4 | 0.92 | 0.75 | $2.43 \times 10^{-9}$ | $5.72 \times 10^{-10}$ |
| 100 | 200 | 15 | 16.6 | 1.26 | 1.02 | $1.18 \times 10^{-10}$ | $2.78 \times 10^{-12}$ |
| 150 | 300 | 13 | 15.8 | 2.05 | 2.33 | $9.53 \times 10^{-9}$ | $6.57 \times 10^{-11}$ |
| 200 | 400 | 12 | 13.2 | 4.51 | 3.94 | $5.82 \times 10^{-9}$ | $3.44 \times 10^{-10}$ |

Table 7. Numerical results of Algorithm 3.1 for the problem (5.3).

| $m$ | $n$ | mIT | AIT | MCPU | ACPU | AFV | mFV |
| ---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 50 | 100 | 13 | 14.2 | 1.21 | 0.81 | $5.94 \times 10^{-9}$ | $6.38 \times 10^{-10}$ |
| 100 | 200 | 15 | 18.4 | 3.77 | 1.65 | $4.78 \times 10^{-9}$ | $3.94 \times 10^{-11}$ |
| 150 | 300 | 14 | 17.2 | 5.18 | 3.63 | $1.71 \times 10^{-9}$ | $6.79 \times 10^{-11}$ |
| 200 | 400 | 16 | 19.6 | 5.96 | 4.88 | $1.32 \times 10^{-10}$ | $3.02 \times 10^{-10}$ |

Table 8. Numerical results of SDPT3 for the problem (5.3).

## 6. Conclusions

It has been shown in [33] that the Qi-Sun-Zhou smoothing Newton method [27] performs very efficiently for solving complementarity problems in practice. In this paper, by modifying and extending the Qi-Sun-Zhou method [27], we propose a smoothing Newton method for solving the SOCCP based on a new smoothing function. This new function is coercive under suitable conditions, which plays an important role in the convergence analysis. We prove that the proposed algorithm is globally and locally quadratically convergent under suitable assumptions. Some numerical results are also reported which indicate that our algorithm is effective for solving the SOCCP.

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