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# ON CIRCUMCENTER MAPPINGS INDUCED BY NONEXPANSIVE OPERATORS 

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#### Abstract

We introduce the circumcenter mapping induced by a set of (usually nonexpansive) operators. One prominent example of a circumcenter mapping is the celebrated Douglas-Rachford splitting operator. Our study is motivated by the Circumcentered-Douglas-Rachford method recently introduced by Behling, Bello Cruz, and Santos in order to accelerate the Douglas-Rachford method for solving certain classes of feasibility problems. We systematically explore the properness of the circumcenter mapping induced by reflectors or projectors. Numerous examples are presented. We also present a version of Browder's demiclosedness principle for circumcenter mappings.


## 1. Introduction

Throughout this paper, we assume that

$$
\mathcal{H} \text { is a real Hilbert space }
$$

with inner product $\langle\cdot, \cdot\rangle$ and induced norm $\|\cdot\|$. Let $m \in \mathbb{N} \backslash\{0\}$, and let $T_{1}, \ldots, T_{m-1}, T_{m}$ be operators from $\mathcal{H}$ to $\mathcal{H}$. Set

$$
\mathcal{S}=\left\{T_{1}, \ldots, T_{m-1}, T_{m}\right\}
$$

and denote the power set of $\mathcal{H}$ as $2^{\mathcal{H}}$. The associated set-valued operator $\mathcal{S}: \mathcal{H} \rightarrow 2^{\mathcal{H}}$ is defined by

$$
(\forall x \in \mathcal{H}) \quad \mathcal{S}(x)=\left\{T_{1} x, \ldots, T_{m-1} x, T_{m} x\right\}
$$

Unless otherwise specified, we assume that

$$
U_{1}, \ldots, U_{m} \text { are closed affine subspaces of } \mathcal{H}, \text { with } \bigcap_{i=1}^{m} U_{i} \neq \varnothing \text {. }
$$

In this paper, we introduce the circumcenter mapping $C C_{\mathcal{S}}$ induced by $\mathcal{S}$ which maps every element $x \in \mathcal{H}$ to either empty set or the (unique if it exists) circumcenter of the finitely many elements in the nonempty set $\mathcal{S}(x)$. In fact, the circumcenter mapping $C C_{\mathcal{S}}$ induced by $\mathcal{S}$ is the composition $C C \circ \mathcal{S}$ where $C C$ is the circumcenter operator defined in [4]. The domain of $C C_{\mathcal{S}}$ is defined to be $\operatorname{dom} C C_{\mathcal{S}}=\left\{x \in \mathcal{H} \mid C C_{\mathcal{S}} x \neq \varnothing\right\}$. We say the circumcenter mapping $C C_{\mathcal{S}}$ is proper, if dom $C C_{\mathcal{S}}=\mathcal{H}$. Properness is an important property for algorithms where one wishes to consider sequences of the form $\left(C C_{\mathcal{S}}^{k} x\right)_{k \in \mathbb{N}}$.

[^0]The goal of this paper is to explore conditions sufficient for the circumcenter mapping to be proper. We also connect the circumcenter mapping to the celebrated demiclosedness principle by Felix Browder.

The CRM (Circumcentered-Reflection Method) operator $C$ recently investigated by Behling, Bello Cruz, and Santos in [7, page 159] is a particular instance of a proper circumcenter mapping. The C-DRM (Circumcentered-Douglas-Rachford Method) operator $C_{T}$ defined by Behling et al. in [6, Section 2] is CRM operator associated with only two linear subspaces. Hence, the $C_{T}$ is a special case of our proper circumcenter mapping as well.

Behling et al. introduced in [6] the C-DRM which generates iterates by taking the intersection of bisectors of reflection steps to accelerate the Douglas-Rachford method to solve certain classes of feasibility problems. Our paper [4] and this paper are motivated by [6]. The proof of one of our main results, Theorem 4.3, is inspired by [6, Lemma 2]. We now discuss further results and the organization of this paper. In Section 2, we collect various results for subsequent use. In particular, facts on circumcenter operator defined in [4, Definition 3.4] are reviewed in Section 2.3. In Section 3, we introduce the circumcenter mapping $C C_{\mathcal{S}}$ induced by a set of operators $\mathcal{S}$. Based on some known results of circumcenter operator, we derive some sufficient conditions for the circumcenter mapping $C C_{\mathcal{S}}$ to be proper. When $\mathcal{S}$ consists of only three operators, we provide a sufficient and necessary condition for the $C C_{\mathcal{S}}$ to be proper. We also obtain conditions sufficient for continuity. Examples illustrating the tightness of our assumptions are provided as well. Section 3.4 contains the demiclosedness principle for certain circumcenter mappings. In Section 4, we consider the circumcenter of finite subsets drawn from the affine hull of compositions of reflectors. Inspired by [6, Lemma 2], we prove the properness of a certain class of circumcenter mappings induced by reflectors. We also provide improper examples. Two particular instances of $C C_{\mathcal{S}}$, one of which belongs to the class of C-DRM operators from [6] while the other is new, are considered. Comparing to the Douglas-Rachford Method (DRM) and the Method of Alternating projections (MAP), we find in preliminary numerical explorations that $\left(C C_{\mathcal{S}}^{k} x\right)_{k \in \mathbb{N}}$ can be used to solve best approximation problems. It is interesting that in general $C C_{\mathcal{S}}$ is neither continuous nor linear. In Section 5, the operators in $\mathcal{S}$ are chosen from the affine hull of the set of compositions of projectors. We provide both proper and improper examples of corresponding circumcenter mappings. The final Section 6 deals with reflectors and reflected resolvents.

Let us turn to notation. Let $K$ and $C$ be subsets of $\mathcal{H}, z \in \mathcal{H}$ and $\lambda \in \mathbb{R}$. Then $K+C=\{x+y \mid x \in K, y \in C\}, K+z=K+\{z\}$, and $\lambda K=\{\lambda x \mid x \in K\}$. The cardinality of the set $K$ is denoted as $\operatorname{card}(K)$. The intersection of all the linear subspaces of $\mathcal{H}$ containing $K$ is called the span of $K$, and is denoted by span $K$. A nonempty subset $K$ of $\mathcal{H}$ is an affine subspace of $\mathcal{H}$ if $(\forall \rho \in \mathbb{R}) \rho K+(1-\rho) K=K$; moreover, the smallest affine subspace containing $K$ is the affine hull of $K$, denoted aff $K$. Assume that $C$ is a nonempty closed, convex subset in $\mathcal{H}$. We denote by $\mathrm{P}_{C}$ the projector onto $C . \mathrm{R}_{C}:=2 \mathrm{P}_{C}-\mathrm{Id}$ is the reflector associated with $C$. Let $T: \mathcal{H} \rightarrow \mathcal{H}$. The set of fixed points of $T$ is $\operatorname{Fix} T=\{x \in \mathcal{H} \mid x=T x\}$. Let $\left(x_{k}\right)_{k \in \mathbb{N}}$ be a sequence in $\mathcal{H}$ and let $x \in \mathcal{H}$. We use $x_{k} \rightharpoonup x$ to indicate that $\left(x_{k}\right)_{k \in \mathbb{N}}$ converges weakly to $x$. The set $\mathbf{B}[x ; r]:=\{y \in \mathcal{H} \mid\|y-x\| \leq r\}$ is the closed ball
centered at $x$ of radius $r \geq 0$. For other notation not explicitly defined here, we refer the reader to [2].

## 2. Auxiliary Results

In this section, we provide various results that will be useful in the sequel. We start with some facts about affine subspaces.

### 2.1. Affine subspaces and related concepts.

Definition 2.1 ([10, page 4]). An affine subspace $C$ is said to be parallel to an affine subspace $M$ if $C=M+a$ for some $a \in \mathcal{H}$.

Fact 2.2 ([10, Theorem 1.2]). Every affine subspace $C$ is parallel to a unique linear subspace $L$, which is given by

$$
(\forall y \in C) \quad L=C-y=C-C .
$$

Definition 2.3 ([10, page 4]). The dimension of a nonempty affine subspace is defined to be the dimension of the linear subspace parallel to it.

Fact $2.4([10$, page 7$])$. Let $x_{1}, \ldots, x_{m} \in \mathcal{H}$. Then the affine hull is given by

$$
\operatorname{aff}\left\{x_{1}, \ldots, x_{m}\right\}=\left\{\lambda_{1} x_{1}+\cdots+\lambda_{m} x_{m} \mid \lambda_{1}, \ldots, \lambda_{m} \in \mathbb{R} \text { and } \sum_{i=1}^{m} \lambda_{i}=1\right\}
$$

Fact 2.5 ([4, Lemma 2.6]). Let $x_{1}, \ldots, x_{m} \in \mathcal{H}$, where $m \geq 2$. Then for every $i_{0} \in\{2, \ldots, m\}$, we have

$$
\begin{aligned}
& \operatorname{aff}\left\{x_{1}, \ldots, x_{m}\right\} \\
= & x_{1}+\operatorname{span}\left\{x_{2}-x_{1}, \ldots, x_{m}-x_{1}\right\} \\
= & x_{i_{0}}+\operatorname{span}\left\{x_{1}-x_{i_{0}}, \ldots, x_{i_{0}-1}-x_{i_{0}}, x_{i_{0}+1}-x_{i_{0}}, \ldots, x_{m}-x_{i_{0}}\right\} .
\end{aligned}
$$

Definition 2.6 ([10, page 6]). Let $x_{0}, x_{1}, \ldots, x_{m} \in \mathcal{H}$. The $m+1$ vectors $x_{0}, x_{1}, \ldots$, $x_{m}$ are said to be affinely independent if aff $\left\{x_{0}, x_{1}, \ldots, x_{m}\right\}$ is $m$-dimensional. We will also say $\left(x_{0}, x_{1}, \ldots, x_{m}\right)=\left(x_{i}\right)_{i \in\{0,1, \ldots, m\}}$ is affinely independent.
Fact 2.7 ([10, page 7$])$. Let $x_{1}, x_{2}, \ldots, x_{m} \in \mathcal{H}$. Then $x_{1}, x_{2}, \ldots, x_{m}$ are affinely independent if and only if $x_{2}-x_{1}, \ldots, x_{m}-x_{1}$ are linearly independent.
2.2. Projectors and reflectors. Our first result follows easily from the definitions.

Lemma 2.8. Let $C$ be a nonempty closed convex subset of $\mathcal{H}$. Then
(i) $\mathrm{P}_{C} \mathrm{P}_{C}=\mathrm{P}_{C}$.
(ii) If $C$ is a closed affine subspace, then $\mathrm{R}_{C} \mathrm{R}_{C}=\mathrm{Id}$.

Fact 2.9 ([9, Theorem 5.8]). Let $C$ be a closed linear subspace of $\mathcal{H}$. Then
(i) $\mathrm{Id}=\mathrm{P}_{C}+\mathrm{P}_{C^{\perp}}$.
(ii) $C^{\perp}=\left\{x \in \mathcal{H} \mid \mathrm{P}_{C}(x)=0\right\}$ and $C=\left\{x \in \mathcal{H} \mid \mathrm{P}_{C^{\perp}}(x)=0\right\}=\{x \in$ $\left.\mathcal{H} \mid \mathrm{P}_{C}(x)=x\right\}$.

The following result is a mild extension [6, Proposition 1] and it is useful in the proof of Theorem 4.3.

Proposition 2.10. Let $C$ be a closed affine subspace of $\mathcal{H}$. Then the following hold:
(i) The projector $\mathrm{P}_{C}$ and the reflector $\mathrm{R}_{C}$ are affine operators.
(ii) Let $x$ be in $\mathcal{H}$ and let $p$ be in $\mathcal{H}$. Then

$$
p=\mathrm{P}_{C} x \Longleftrightarrow p \in C \quad \text { and } \quad(\forall v \in C)(\forall w \in C) \quad\langle x-p, v-w\rangle=0
$$

(iii) $(\forall x \in \mathcal{H})(\forall v \in C)\left\|x-\mathrm{P}_{C} x\right\|^{2}+\left\|v-\mathrm{P}_{C} x\right\|^{2}=\|x-v\|^{2}$.
(iv) $(\forall x \in \mathcal{H})(\forall y \in \mathcal{H})\|x-y\|=\left\|\mathrm{R}_{C} x-\mathrm{R}_{C} y\right\|$.
(v) $(\forall x \in \mathcal{H})(\forall v \in C)\|x-v\|=\left\|\mathrm{R}_{C} x-v\right\|$.

Proof. (i): $\mathrm{P}_{C}$ is affine by [2, Corollary $\left.3.22(\mathrm{ii})\right]$; this implies that $\mathrm{R}_{C}=2 \mathrm{P}_{C}-\mathrm{Id}$ is affine as well.
(ii): [2, Corollary 3.22(i)].
(iii): Indeed, for every $x \in \mathcal{H}$ and $v \in C$,

$$
\begin{aligned}
\|x-v\|^{2} & =\left\|x-\mathrm{P}_{C} x-\left(v-\mathrm{P}_{C} x\right)\right\|^{2} \\
& =\left\|x-\mathrm{P}_{C} x\right\|^{2}-2\left\langle x-\mathrm{P}_{C} x, v-\mathrm{P}_{C} x\right\rangle+\left\|v-\mathrm{P}_{C} x\right\|^{2} \\
& =\left\|x-\mathrm{P}_{C} x\right\|^{2}+\left\|v-\mathrm{P}_{C} x\right\|^{2} . \quad(\text { by }(\mathrm{ii}))
\end{aligned}
$$

(iv): For every $x \in \mathcal{H}$, and for every $y \in \mathcal{H}$, by (ii),

$$
\begin{aligned}
& \left\langle\mathrm{P}_{C} x-\mathrm{P}_{C} y, \mathrm{P}_{C} x-x\right\rangle-\left\langle\mathrm{P}_{C} x-\mathrm{P}_{C} y, \mathrm{P}_{C} y-y\right\rangle=0 \\
\Longleftrightarrow & \left\langle\mathrm{P}_{C} x-\mathrm{P}_{C} y, \mathrm{P}_{C} x-\mathrm{P}_{C} y-(x-y)\right\rangle=0 \\
\Longleftrightarrow & \|x-y\|^{2}=4\left\|\mathrm{P}_{C} x-\mathrm{P}_{C} y\right\|^{2}-4\left\langle\mathrm{P}_{C} x-\mathrm{P}_{C} y, x-y\right\rangle+\|x-y\|^{2} \\
\Longleftrightarrow & \|x-y\|^{2}=\left\|\left(2 \mathrm{P}_{C} x-x\right)-\left(2 \mathrm{P}_{C} y-y\right)\right\|^{2} \\
\Longleftrightarrow & \|x-y\|=\left\|\mathrm{R}_{C} x-\mathrm{R}_{C} y\right\| . \quad\left(\text { by } \mathrm{R}_{C}=2 \mathrm{P}_{C}-\mathrm{Id}\right)
\end{aligned}
$$

(v): Notice that Fix $\mathrm{R}_{C}=C$ and then use (iv).
2.3. Circumcenters. In this subsection, $\mathcal{P}(\mathcal{H})$ is the set of all nonempty subsets of $\mathcal{H}$ containing finitely many elements.

By [4, Proposition 3.3], we know that for every $K \in \mathcal{P}(\mathcal{H})$, there is at most one point $p \in \operatorname{aff}(K)$ such that $\{\|p-x\| \mid x \in K\}$ is a singleton. Hence, the following notion is well-defined.

Definition 2.11 (circumcenter operator ([4, Definition 3.4])). The circumcenter operator is $C C: \mathcal{P}(\mathcal{H}) \rightarrow \mathcal{H} \cup\{\varnothing\}$ such that for every $K \in \mathcal{P}(\mathcal{H})$,

$$
C C(K)=\left\{\begin{array}{l}
p, \text { if } p \in \operatorname{aff}(K) \text { and }\{\|p-x\| \mid x \in K\} \text { is a singleton; } \\
\varnothing, \text { otherwise }
\end{array}\right.
$$

In particular, when $C C(K) \in \mathcal{H}$, that is, $C C(K) \neq \varnothing$, we say that the circumcenter of $K$ exists and we call $C C(K)$ the circumcenter of $K$.

Fact 2.12 ([4, Example 3.6]). Let $x_{1}, x_{2}$ be in $\mathcal{H}$. Then

$$
C C\left(\left\{x_{1}, x_{2}\right\}\right)=\frac{x_{1}+x_{2}}{2}
$$

Fact $2.13\left(\left[4\right.\right.$, Theorem 4.1]). Set $K=\left\{x_{1}, \ldots, x_{m}\right\}$, where $x_{1}, \ldots, x_{m}$ are affinely independent. Then $C C(K) \in \mathcal{H}$, which means that $C C(K)$ is the unique point satisfying the following two conditions:
(i) $C C(K) \in \operatorname{aff}(K)$, and
(ii) $\{\|C C(K)-y\| \mid y \in K\}$ is a singleton.

Moreover,

$$
\begin{aligned}
& C C(K) \\
= & x_{1}+\frac{1}{2}\left(x_{2}-x_{1}, \ldots, x_{m}-x_{1}\right) G\left(x_{2}-x_{1}, \ldots, x_{m}-x_{1}\right)^{-1}\left(\begin{array}{c}
\left\|x_{2}-x_{1}\right\|^{2} \\
\vdots \\
\left\|x_{m}-x_{1}\right\|^{2}
\end{array}\right)
\end{aligned}
$$

where $G\left(x_{2}-x_{1}, \ldots, x_{m-1}-x_{1}, x_{m}-x_{1}\right)$ is the Gram matrix of $x_{2}-x_{1}, \ldots, x_{m-1}-$ $x_{1}, x_{m}-x_{1}$, i.e.,

$$
\begin{aligned}
& G\left(x_{2}-x_{1}, \ldots, x_{m-1}-x_{1}, x_{m}-x_{1}\right) \\
= & \left(\begin{array}{cccc}
\left\|x_{2}-x_{1}\right\|^{2} & \left\langle x_{2}-x_{1}, x_{3}-x_{1}\right\rangle & \cdots & \left\langle x_{2}-x_{1}, x_{m}-x_{1}\right\rangle \\
\vdots & \vdots & & \vdots \\
\left\langle x_{m-1}-x_{1}, x_{2}-x_{1}\right\rangle & \left\langle x_{m-1}-x_{1}, x_{3}-x_{1}\right\rangle & \cdots & \left\langle x_{m-1}-x_{1}, x_{m}-x_{1}\right\rangle \\
\left\langle x_{m}-x_{1}, x_{2}-x_{1}\right\rangle & \left\langle x_{m}-x_{1}, x_{3}-x_{1}\right\rangle & \cdots & \left\|x_{m}-x_{1}\right\|^{2}
\end{array}\right) .
\end{aligned}
$$

Fact 2.14 ([4, Theorem 8.1]). Suppose that $K=\{x, y, z\} \in \mathcal{P}(\mathcal{H})$ and that $\operatorname{card}(K)=3$. Then $x, y, z$ are affinely independent if and only if $C C(K) \in \mathcal{H}$.

Combining Fact 2.12 and Fact 2.14, we obtain the following two results.
Corollary 2.15. Let $K=\left\{x_{1}, x_{2}, x_{3}\right\} \in \mathcal{P}(\mathcal{H})$. Then $C C(K) \in \mathcal{H}$ if and only if exactly one of the following cases holds.
(i) $\operatorname{card}\left\{x_{1}, x_{2}, x_{3}\right\} \leq 2$.
(ii) card $\left\{x_{1}, x_{2}, x_{3}\right\}=3$ and if there is $\{\alpha, \beta\} \subseteq \mathbb{R}$ such that $\alpha\left(x_{2}-x_{1}\right)+\beta\left(x_{3}-\right.$ $\left.x_{1}\right)=0$, then $\alpha=0$ and $\beta=0$.

Corollary 2.16. Let $a, b, c$ be in $\mathbb{R}$. Then there exists no $x \in \mathbb{R}$ such that $|x-a|=$ $|x-b|=|x-c|$ if and only if $\operatorname{card}\{a, b, c\}=3$.

Fact 2.17 (scalar multiples $([4$, Proposition 6.1])). Let $K \in \mathcal{P}(\mathcal{H})$ and $\lambda \in \mathbb{R} \backslash\{0\}$. Then $C C(\lambda K)=\lambda C C(K)$.

Fact 2.18 (translations $([4$, Proposition 6.3])). Let $K \in \mathcal{P}(\mathcal{H})$ and $y \in \mathcal{H}$. Then $C C(K+y)=C C(K)+y$.
Fact 2.19 ([4, Lemma 4.2]). Let $K \in \mathcal{P}(\mathcal{H})$ and let $M \subseteq K$ be such that $\operatorname{aff}(M)=$ $\operatorname{aff}(K)$. Suppose that $C C(K) \in \mathcal{H}$. Then $C C(K)=C C(M)$.

Fact 2.20 ([4, Theorem 7.1]). Let $K=\left\{x_{1}, \ldots, x_{m}\right\} \in \mathcal{P}(\mathcal{H})$. Suppose that $C C(K) \in \mathcal{H}$. Then the following hold.
(i) Let $t=\operatorname{dim}\left(\operatorname{span}\left\{x_{2}-x_{1}, \ldots, x_{m}-x_{1}\right\}\right), \widetilde{K}=\left\{x_{1}, x_{i_{1}}, \ldots, x_{i_{t}}\right\} \subseteq K$ be such that $x_{i_{1}}-x_{1}, \ldots, x_{i_{t}}-x_{1}$ form a basis of $\operatorname{span}\left\{x_{2}-x_{1}, \ldots, x_{m}-x_{1}\right\}$.

Furthermore, let $\left(\left(x_{1}^{(k)}, x_{i_{1}}^{(k)}, \ldots, x_{i_{t}}^{(k)}\right)\right)_{k \geq 1} \subseteq \mathcal{H}^{t+1}$ with

$$
\lim _{k \rightarrow \infty}\left(x_{1}^{(k)}, x_{i_{1}}^{(k)}, \ldots, x_{i_{t}}^{(k)}\right)=\left(x_{1}, x_{i_{1}}, \ldots, x_{i_{t}}\right)
$$

and set $(\forall k \geq 1) \widetilde{K}^{(k)}=\left\{x_{1}^{(k)}, x_{i_{1}}^{(k)}, \ldots, x_{i_{t}}^{(k)}\right\}$. Then there exist $N \in \mathbb{N}$ such that for every $k \geq N, C C\left(\widetilde{K}^{(k)}\right) \in \mathcal{H}$ and

$$
\lim _{k \rightarrow \infty} C C\left(\widetilde{K}^{(k)}\right)=C C(\widetilde{K})=C C(K)
$$

(ii) Suppose that $x_{1}, \ldots, x_{m-1}, x_{m}$ are affinely independent, and let

$$
\left(\left(x_{1}^{(k)}, \ldots, x_{m-1}^{(k)}, x_{m}^{(k)}\right)\right)_{k \geq 1} \subseteq \mathcal{H}^{m}
$$

satisfy $\lim _{k \rightarrow \infty}\left(x_{1}^{(k)}, \ldots, x_{m-1}^{(k)}, x_{m}^{(k)}\right)=\left(x_{1}, \ldots, x_{m-1}, x_{m}\right)$. Assume that $(\forall k \geq 1) K^{(k)}=\left\{x_{1}^{(k)}, \ldots, x_{m-1}^{(k)}, x_{m}^{(k)}\right\}$. Then

$$
\lim _{k \rightarrow \infty} C C\left(K^{(k)}\right)=C C(K)
$$

Fact 2.21 ([4, Example 7.6]). Suppose that $\mathcal{H}=\mathbb{R}^{2}$. Let $x_{1}=(-2,0)$ and $x_{2}=$ $x_{3}=(2,0)$. Let $(\forall k \geq 1)\left(x_{1}^{(k)}, x_{2}^{(k)}, x_{3}^{(k)}\right)=\left((-2,0),(2,0),\left(2-\frac{1}{k}, \frac{1}{4 k}\right)\right)$. Then

$$
(\forall k \geq 1) \quad C C\left(\left\{x_{1}^{(k)}, x_{2}^{(k)}, x_{3}^{(k)}\right\}\right)=\left(0,-8+\frac{2}{k}+\frac{1}{8 k}\right)
$$

## 3. Circumcenter mappings induced by operators

Suppose that $T_{1}, \ldots, T_{m-1}, T_{m}$ are operators from $\mathcal{H}$ to $\mathcal{H}$, with $m \in \mathbb{N} \backslash\{0\}$ and that

$$
\mathcal{S}=\left\{T_{1}, \ldots, T_{m-1}, T_{m}\right\} \text { and }(\forall x \in \mathcal{H}) \mathcal{S}(x)=\left\{T_{1} x, \ldots, T_{m-1} x, T_{m} x\right\}
$$

### 3.1. Definition.

Definition 3.1 (induced circumcenter mapping). The circumcenter mapping $C C_{\mathcal{S}}$ induced by $\mathcal{S}$ is

$$
C C_{\mathcal{S}}: \mathcal{H} \rightarrow \mathcal{H} \cup\{\varnothing\}: x \mapsto C C(\mathcal{S}(x))
$$

that is, $C C_{\mathcal{S}}=C C \circ \mathcal{S}$. The domain of $C C_{\mathcal{S}}$ is

$$
\operatorname{dom} C C_{\mathcal{S}}=\left\{x \in \mathcal{H} \mid C C_{\mathcal{S}} x \neq \varnothing\right\}
$$

Recall that if dom $C C_{\mathcal{S}}=\mathcal{H}$, then we say the circumcenter mapping $C C_{\mathcal{S}}$ induced by $\mathcal{S}$ is proper; otherwise, we will call $C C_{\mathcal{S}}$ improper.

Remark 3.2. By Definitions 2.11 and 3.1 , for every $x \in \mathcal{H}$, if the circumcenter of the set $\mathcal{S}(x)$ defined in Definition 2.11 does not exist in $\mathcal{H}$, then $C C_{\mathcal{S}} x=\varnothing$. Otherwise, $C C_{\mathcal{S}} x$ is the unique point satisfying the two conditions below:
(i) $C C_{\mathcal{S}} x \in \operatorname{aff}(\mathcal{S}(x))=\operatorname{aff}\left\{T_{1} x, \ldots, T_{m-1} x, T_{m} x\right\}$, and
(ii) $\left\|C C_{\mathcal{S}} x-T_{1} x\right\|=\cdots=\left\|C C_{\mathcal{S}} x-T_{m-1} x\right\|=\left\|C C_{\mathcal{S}} x-T_{m} x\right\|$.
3.2. Basic properties. We start with some examples.

Proposition 3.3. Assume $\mathcal{S}=\left\{T_{1}, T_{2}\right\}$. Then $C C_{\mathcal{S}}$ is proper. Moreover,

$$
(\forall x \in \mathcal{H}) \quad C C_{\mathcal{S}} x=\frac{T_{1} x+T_{2} x}{2}
$$

Proof. Clear from Fact 2.12 and Definition 3.1.
Corollary 3.4. Let $\mathcal{S}=\left\{T_{1}, T_{2}, T_{3}\right\}$ and let $x \in \mathcal{H}$. Then $x \notin \operatorname{dom} C C_{\mathcal{S}}$ if and only if $\operatorname{card}\left\{T_{1} x, T_{2} x, T_{3} x\right\}=3$ and there exists $(\alpha, \beta) \in \mathbb{R}^{2} \backslash\{(0,0)\}$ such that $\alpha\left(T_{2} x-T_{1} x\right)+\beta\left(T_{3} x-T_{1} x\right)=0$.

Proof. This follows from Corollary 2.15.
Example 3.5. Assume that $\mathcal{H}=\mathbb{R}^{2}$. Set $U_{1}=\mathbb{R} \cdot(1,0), U_{2}=\mathbb{R} \cdot(0,1)$, and let $\alpha \in \mathbb{R}$. Set $\mathcal{S}=\left\{\alpha \mathrm{Id}, R_{U_{1}}, R_{U_{2}}\right\}$. Then the following hold:
(i) If $\alpha=0$, then $\operatorname{dom} C C_{\mathcal{S}}=\{(0,0)\}$.
(ii) If $\alpha=1$ or $\alpha=-1$, then $\operatorname{dom} C C_{\mathcal{S}}=\mathbb{R}^{2}$, i.e., $C C_{\mathcal{S}}$ is proper.
(iii) If $\alpha \in \mathbb{R} \backslash\{0,1,-1\}$, then $\operatorname{dom} C C_{\mathcal{S}}=\left(\mathbb{R}^{2} \backslash\left(U_{1} \cup U_{2}\right)\right) \cup\{(0,0)\}$.

Proposition 3.6. Suppose that for every $x \in \mathcal{H}$, there exists a point $p(x) \in \mathcal{H}$ such that
(i) $p(x) \in \operatorname{aff}\left\{T_{1} x, \ldots, T_{m-1} x, T_{m} x\right\}$, and
(ii) $\left\|p(x)-T_{1} x\right\|=\cdots=\left\|p(x)-T_{m-1} x\right\|=\left\|p(x)-T_{m} x\right\|$.

Then $C C_{\mathcal{S}}$ is proper and

$$
(\forall x \in \mathcal{H}) \quad C C_{\mathcal{S}} x=p(x)
$$

Proof. This follows from Remark 3.2.
Proposition 3.7. Suppose that for every $x \in \mathcal{H}$, there exists $\mathrm{I}(x) \subseteq \mathrm{I}:=\{1, \ldots, m\}$ such that $\operatorname{card}(\mathrm{I}(x))=\operatorname{card}(\mathcal{S}(x))$ and $\left(T_{i} x\right)_{i \in \mathrm{I}(x)}$ is affinely independent. Then $C C_{\mathcal{S}}$ is proper.
Proof. Let $x \in \mathcal{H}$. Since $\mathrm{I}(x) \subseteq \mathrm{I}$, we have $\left\{T_{i} x\right\}_{i \in \mathrm{I}(x)} \subseteq \mathcal{S}(x)$. The affine independence of $\left(T_{i} x\right)_{i \in \mathrm{I}(x)}$ yields card $\left(\left\{T_{i} x\right\}_{i \in \mathrm{I}(x)}\right)=\operatorname{card}(\mathrm{I}(x))$. Combining with card $(\mathrm{I}(x))=\operatorname{card}(\mathcal{S}(x))$, we obtain that $\left\{T_{i} x\right\}_{i \in \mathrm{I}(x)}=\mathcal{S}(x)$, which implies that

$$
\begin{equation*}
C C_{\mathcal{S}} x=C C(\mathcal{S}(x))=C C\left(\left\{T_{i} x\right\}_{i \in \mathrm{I}(x)}\right) \tag{3.1}
\end{equation*}
$$

Using the assumption that $\left(T_{i} x\right)_{i \in \mathrm{I}(x)}$ is affinely independent, by Fact 2.13 , we deduce that $C C\left(\left\{T_{i} x\right\}_{i \in \mathrm{I}(x)}\right) \in \mathcal{H}$. Combining with (3.1), we deduce that $(\forall x \in \mathcal{H})$ $C C_{\mathcal{S}} x \in \mathcal{H}$, i.e., $C C_{\mathcal{S}}$ is proper.

The following example illustrates that the converse of Proposition 3.7 is not true in general.
Example 3.8. Let $U$ be a closed linear subspace of $\mathcal{H}$ with $\{0\} \neq U \varsubsetneqq \mathcal{H}$. Denote by 0 also the zero operator: $(\forall x \in \mathcal{H}) 0(x)=0$. Set $\mathcal{S}=\left\{\operatorname{Id}, \mathrm{P}_{U}, \mathrm{P}_{U^{\perp}}, 0\right\}$. Then the following hold:
(i) $(\forall x \in \mathcal{H}) C C_{\mathcal{S}} x=\frac{x}{2}$; consequently, $C C_{\mathcal{S}}$ is proper.
(ii) $\left(\forall x \in \mathcal{H} \backslash\left(U \cup U^{\perp}\right)\right) x, \mathrm{P}_{U} x, \mathrm{P}_{U^{\perp}} x, 0(x)$ are pairwise distinct.
(iii) $(\forall x \in \mathcal{H}) \mathrm{Id} x, \mathrm{P}_{U} x, \mathrm{P}_{U^{\perp}} x, 0(x)$ are affinely dependent.

Proof. (i): Let $x \in \mathcal{H}$. By Proposition 2.10 (iii) and by $0 \in U$ and $\mathrm{P}_{U} x \in U$, we deduce that $\left\|\frac{x}{2}-\mathrm{P}_{U} \frac{x}{2}\right\|^{2}+\left\|\mathrm{P}_{U} \frac{x}{2}\right\|^{2}=\left\|\frac{x}{2}\right\|^{2}$ and that $\left\|\frac{x}{2}-\mathrm{P}_{U} \frac{x}{2}\right\|^{2}+\left\|\mathrm{P}_{U} x-\mathrm{P}_{U} \frac{x}{2}\right\|^{2}=$ $\left\|\frac{x}{2}-\mathrm{P}_{U} x\right\|^{2}$. Combining with the linearity of $\mathrm{P}_{U}$, we obtain

$$
\begin{equation*}
\left\|\frac{x}{2}\right\|=\left\|\frac{x}{2}-\mathrm{P}_{U} x\right\| \tag{3.2}
\end{equation*}
$$

Similarly, by Proposition 2.10(iii) again, replace $U$ in the above analysis by $U^{\perp}$ to yield that

$$
\begin{equation*}
\left\|\frac{x}{2}\right\|=\left\|\frac{x}{2}-\mathrm{P}_{U^{\perp}} x\right\| \tag{3.3}
\end{equation*}
$$

Combining (3.2) with (3.3), we obtain that

$$
\begin{equation*}
\left\|\frac{x}{2}\right\|=\left\|\frac{x}{2}-0(x)\right\|=\left\|\frac{x}{2}-x\right\|=\left\|\frac{x}{2}-\mathrm{P}_{U} x\right\|=\left\|\frac{x}{2}-\mathrm{P}_{U^{\perp}} x\right\| \tag{3.4}
\end{equation*}
$$

Since $\frac{x}{2}=\frac{x}{2}+\frac{0}{2} \in \operatorname{aff}\left\{x, \mathrm{P}_{U} x, \mathrm{P}_{U \perp} x, 0(x)\right\},(3.4)$ yields that $(\forall x \in \mathcal{H}) C C_{\mathcal{S}} x=\frac{x}{2}$.
(ii): In fact, by Fact 2.9(ii),

$$
\begin{align*}
& x=\mathrm{P}_{U} x \Longleftrightarrow x \in U  \tag{3.5a}\\
& x=\mathrm{P}_{U^{\perp}} x \Longleftrightarrow x \in U^{\perp}  \tag{3.5b}\\
& U \cap U^{\perp}=\{0\} \tag{3.5c}
\end{align*}
$$

In addition, Combining (3.5) with Fact 2.9(i), we know that

$$
\mathrm{P}_{U} x=\mathrm{P}_{U^{\perp}} x \Longrightarrow \mathrm{P}_{U} x=\mathrm{P}_{U^{\perp}} x=0 \Longrightarrow x=\mathrm{P}_{U} x+\mathrm{P}_{U^{\perp}} x=0 \in U \cup U^{\perp}
$$

Hence, for every $x \in \mathcal{H} \backslash\left(U \cup U^{\perp}\right), x, \mathrm{P}_{U} x, \mathrm{P}_{U^{\perp}} x, 0(x)$ are pairwise distinct.
(iii): Now for every $x \in \mathcal{H}$,

$$
\begin{aligned}
& x=\mathrm{P}_{U} x+\mathrm{P}_{U^{\perp}} x \\
\Rightarrow & x, \mathrm{P}_{U} x, \mathrm{P}_{U^{\perp}} x \text { are linear dependent } \\
\Leftrightarrow & x-0, \mathrm{P}_{U} x-0, \mathrm{P}_{U^{\perp}} x-0 \text { are linear dependent } \\
\Leftrightarrow & \left.0(x), \operatorname{Id} x, \mathrm{P}_{U} x, \mathrm{P}_{U^{\perp}} x \text { are affinely dependent. (by Fact } 2.7\right)
\end{aligned}
$$

The proof is complete.
The following theorem provides a way to verify the properness of $C C_{\mathcal{S}}$ where $\mathcal{S}$ contains three operators.

Theorem 3.9. Suppose that $\mathcal{S}=\left\{T_{1}, T_{2}, T_{3}\right\}$. Then $C C_{\mathcal{S}}$ is proper if and only if for every $x \in \mathcal{H}$ with card $(\mathcal{S}(x))=3$, the vectors $T_{1} x, T_{2} x, T_{3} x$ are affinely independent.

Proof. By Fact 2.14, for every $x \in \mathcal{H}$ with $\operatorname{card}(\mathcal{S}(x))=3$,

$$
\begin{equation*}
C C_{\mathcal{S}} x \in \mathcal{H} \Longleftrightarrow T_{1} x, T_{2} x, T_{3} x \text { are affinely independent. } \tag{3.6}
\end{equation*}
$$

$" \Longrightarrow "$ It follows directly from (3.6).
" ": Assume that for every $x \in \mathcal{H}$ with $\operatorname{card}(\mathcal{S}(x))=3, T_{1} x, T_{2} x, T_{3} x$ are affinely independent in $\mathcal{H}$. Let $x \in \mathcal{H}$. If $\operatorname{card}(\mathcal{S}(x))=3$, by (3.6) and the
assumption, then $C C_{\mathcal{S}} x \in \mathcal{H}$. Assume card $(\mathcal{S}(x)) \leq 2$, by Proposition $3.3, C C_{\mathcal{S}} x \in$ $\mathcal{H}$. Altogether, $(\forall x \in \mathcal{H}), C C_{\mathcal{S}} x \in \mathcal{H}$, which means that $C C_{\mathcal{S}}$ is proper.
Proposition 3.10. Suppose that $\mathcal{S}=\left\{T_{1}, \ldots, T_{m-1}, T_{m}\right\}$. Then the following hold:
(i) $\cap_{i=1}^{m} \operatorname{Fix} T_{i} \subseteq \operatorname{Fix} C C_{\mathcal{S}}$.
(ii) If Fix $C C_{\mathcal{S}} \subseteq \cup_{i=1}^{m} \operatorname{Fix} T_{i}$, then $\operatorname{Fix} C C_{\mathcal{S}}=\cap_{i=1}^{m} \operatorname{Fix} T_{i}$.
(iii) If $T_{1}=\mathrm{Id}$, then $\cap_{i=1}^{m} \operatorname{Fix} T_{i}=\operatorname{Fix} C C_{\mathcal{S}}$.

Proof. (i): Let $x \in \cap_{i=1}^{m} \operatorname{Fix} T_{i}$. Then

$$
\begin{equation*}
(\forall i \in\{1, \ldots, m-1, m\}) \quad T_{i} x=x \tag{3.7}
\end{equation*}
$$

which yields that $\operatorname{aff}\left\{T_{1} x, \ldots, T_{m-1} x, T_{m} x\right\}=\operatorname{aff}\{x\}=\{x\}$. In addition, by (3.7),

$$
\left\|x-T_{1} x\right\|=\cdots=\left\|x-T_{m-1} x\right\|=\left\|x-T_{m} x\right\|=0
$$

Therefore, we obtain that $C C_{\mathcal{S}} x=x$, which means that $x \in \operatorname{Fix} C C_{\mathcal{S}}$. Hence, $\cap_{i=1}^{m} \operatorname{Fix} T_{i} \subseteq \operatorname{Fix} C C_{\mathcal{S}}$.
(ii): Let $x \in \operatorname{Fix} C C_{\mathcal{S}}$. By the assumption, there is $i_{0} \in\{1, \ldots, m\}$ such that

$$
\begin{equation*}
x=T_{i_{0}} x \tag{3.8}
\end{equation*}
$$

Now $x \in \operatorname{Fix} C C_{\mathcal{S}}$, i.e., $x=C C_{\mathcal{S}} x$, implies that

$$
\begin{equation*}
\left\|x-T_{1} x\right\|=\cdots=\left\|x-T_{m-1} x\right\|=\left\|x-T_{m} x\right\| \tag{3.9}
\end{equation*}
$$

Combining (3.9) with (3.8), we obtain that

$$
\left\|x-T_{1} x\right\|=\cdots=\left\|x-T_{m-1} x\right\|=\left\|x-T_{m} x\right\|=0
$$

which means that $x \in \cap_{i=1}^{m} \operatorname{Fix} T_{i}$. Hence, $\operatorname{Fix} C C_{\mathcal{S}} \subseteq \cap_{i=1}^{m} \operatorname{Fix} T_{i}$. Combining with (i), we deduce that Fix $C C_{\mathcal{S}}=\cap_{i=1}^{m} \operatorname{Fix} T_{i}$.
(iii): If $T_{1}=\mathrm{Id}$, then Fix $T_{1}=\mathcal{H}$ and the result follows from (ii).

Example 3.11. Assume that $\mathcal{H}=\mathbb{R}^{2}$. Set $T_{1}=\mathrm{P}_{\mathbf{B}[(-2,0) ; 1]}, T_{2}=\mathrm{P}_{\mathbf{B}[(0,2) ; 1]}$, $T_{3}=\mathrm{P}_{\mathbf{B}[(2,0) ; 1]}$, and $\mathcal{S}=\left\{T_{1}, T_{2}, T_{3}\right\}$. Then $C C_{\mathcal{S}}$ is proper. Moreover, $\varnothing=$ $\cap_{i=1}^{3} \operatorname{Fix} T_{i} \varsubsetneqq \operatorname{Fix} C C_{\mathcal{S}}=\{(0,0)\}$.
Proof. The properness of $C C_{\mathcal{S}}$ follows from Theorem 3.9 while the rest is a consequence of elementary manipulations.

### 3.3. Continuity.

Proposition 3.12. Assume that the elements of $\mathcal{S}=\left\{T_{1}, \ldots, T_{m-1}, T_{m}\right\}$ are continuous operators and that $x \in \operatorname{dom} C C_{\mathcal{S}}$. Then the following hold:
(i) Let $\widetilde{\mathcal{S}}_{x}=\left\{T_{1}, T_{i_{1}}, \ldots, T_{i_{d_{x}}}\right\} \subseteq \mathcal{S}$ be such that $T_{i_{1}} x-T_{1} x, \ldots, T_{i_{d_{x}}} x-T_{1} x$ is a basis of $\operatorname{span}\left\{T_{2} x-T_{1} x, \ldots, T_{m} x-T_{1} x\right\}$. Then for every $\left(x^{(k)}\right)_{k \in \mathbb{N}} \subseteq \mathcal{H}$ satisfying $\lim _{k \rightarrow \infty} x^{(k)}=x$, there exists $N \in \mathbb{N}$ such that for every $k \geq N$, $C C_{\widetilde{\mathcal{S}}_{x}}\left(x^{(k)}\right) \in \mathcal{H}$. Moreover

$$
\begin{equation*}
\lim _{k \rightarrow \infty} C C_{\widetilde{\mathcal{S}}_{x}}\left(x^{(k)}\right)=C C_{\widetilde{\mathcal{S}}_{x}} x=C C_{\mathcal{S}} x \tag{3.10}
\end{equation*}
$$

(ii) If $T_{1} x, \ldots, T_{m-1} x, T_{m} x$ are affinely independent, then $C C_{\mathcal{S}}$ is continuous at $x$.

Proof. (i): Let $\left(x^{(k)}\right)_{k \in \mathbb{N}} \subseteq \mathcal{H}$ satisfying $\lim _{k \rightarrow \infty} x^{(k)}=x$. Now

$$
\begin{aligned}
& \mathcal{S}=\left\{T_{1}, \ldots, T_{m-1}, T_{m}\right\}, \quad \mathcal{S}(x)=\left\{T_{1}(x), \ldots, T_{m-1} x, T_{m} x\right\} \\
& \widetilde{\mathcal{S}}_{x}=\left\{T_{1}, T_{i_{1}}, \ldots, T_{i_{d_{x}}}\right\}, \quad \widetilde{\mathcal{S}}_{x}(x)=\left\{T_{1} x, T_{i_{1}} x, \ldots, T_{i_{d_{x}}} x\right\} \\
& \widetilde{\mathcal{S}}_{x}\left(x^{(k)}\right)=\left\{T_{1} x^{(k)}, T_{i_{1}} x^{(k)}, \ldots, T_{i_{d_{x}}} x^{(k)}\right\}
\end{aligned}
$$

By Definition 3.1, $C C_{\mathcal{S}} x \in \mathcal{H}$ means that $C C(\mathcal{S}(x)) \in \mathcal{H}$. By assumptions,

$$
T_{i_{1}} x-T_{1} x, \ldots, T_{i_{d_{x}}} x-T_{1} x \text { is a basis of } \operatorname{span}\left\{T_{2} x-T_{1} x, \ldots, T_{m} x-T_{1} x\right\}
$$

Substituting the $K, \widetilde{K}$ and $\widetilde{K}^{(k)}$ in Fact $2.20(\mathrm{i})$ by the above $\mathcal{S}(x), \widetilde{\mathcal{S}}_{x}(x)$ and $\widetilde{\mathcal{S}}_{x}\left(x^{(k)}\right)$ respectively, we obtain the desired results.
(ii): This follows easily from (i).

The next result summarizes conditions under which the proper circumcenter mapping $C C_{\mathcal{S}}$ is continuous at a point $x$.

Proposition 3.13. Assume that the elements of $\mathcal{S}=\left\{T_{1}, \ldots, T_{m-1}, T_{m}\right\}$ are continuous operators and that $C C_{\mathcal{S}}$ is proper. Let $x \in \mathcal{H}$. The following assertions hold:
(i) If $T_{1} x, \ldots, T_{m-1} x, T_{m} x$ are affinely independent, then $C C_{\mathcal{S}}$ is continuous at $x$.
(ii) If $T_{1} x, \ldots, T_{m-1} x, T_{m} x$ are affinely dependent and $m \leq 2$, then $C C_{\mathcal{S}}$ is continuous at $x$.

Proof. (i) follows from Proposition 3.12 (ii) while (ii) is a consequence of Proposition 3.3.

The following examples show that even when $T_{1} x, \ldots, T_{m-1} x, T_{m} x$ are affinely dependent and $m \geq 3$, then $C C_{S}$ may still be continuous at $x$.
Example 3.14. Suppose that $U$ is a closed linear subspace of $\mathcal{H}$ such that $\{0\} \varsubsetneqq$ $U \varsubsetneqq \mathcal{H}$. Set $\mathcal{S}=\left\{\mathrm{Id}, \mathrm{R}_{U}, \mathrm{R}_{U^{\perp}}\right\}$. Then the following hold:
(i) The vectors $x, \mathrm{R}_{U} x, \mathrm{R}_{U^{\perp}} x$ are affinely dependent for every $x \in U \cup U^{\perp}$.
(ii) $C C_{S} \equiv 0$ which is thus proper and continuous on $\mathcal{H}$.

Proof. (i): For every $x \in U$ (respectively $x \in U^{\perp}$ ), $\mathrm{R}_{U} x=x$ (respectively $\mathrm{R}_{U \perp} x=$ $x$ ), which implies that $x, \mathrm{R}_{U} x, \mathrm{R}_{U^{\perp}} x$, which is $x, x, \mathrm{R}_{U \perp} x$ (respectively $x, \mathrm{R}_{U} x, x$ ) are affinely dependent.
(ii): Since $\mathrm{Id}=\mathrm{P}_{U}+\mathrm{P}_{U^{\perp}}$ and $\mathrm{R}_{U}=2 \mathrm{P}_{U}-\mathrm{Id}$, we have

$$
\begin{aligned}
\frac{\mathrm{R}_{U}+\mathrm{R}_{U^{\perp}}}{2} & =\frac{\left(2 \mathrm{P}_{U}-\mathrm{Id}\right)+\left(2 \mathrm{P}_{U^{\perp}}-\mathrm{Id}\right)}{2} \\
& =\frac{1}{2}\left(2 \mathrm{P}_{U}-\mathrm{Id}+2\left(\mathrm{Id}-\mathrm{P}_{U}\right)-\mathrm{Id}\right)=0
\end{aligned}
$$

Let $x \in \mathcal{H}$. Then $0=\frac{\mathrm{R}_{U} x+\mathrm{R}_{U^{\perp}} x}{2} \in \operatorname{aff}\left\{x, \mathrm{R}_{U} x, \mathrm{R}_{U \perp} x\right\}$. In addition, clearly $0 \in U \cap U^{\perp}$. In Proposition $2.10(\mathrm{v})$, substitute $C=U$, and let the point $v=0$. We

[^1]get $\|x\|=\left\|\mathrm{R}_{U} x\right\|$. Similarly, in Proposition $2.10(\mathrm{v})$, substitute $C=U^{\perp}$ and let the point $v=0$. We get $\|x\|=\left\|\mathrm{R}_{U^{\perp}} x\right\|$. Hence, we have
(a) $0 \in \operatorname{aff}\left\{x, \mathrm{R}_{U} x, \mathrm{R}_{U^{\perp}} x\right\}$ and
(b) $\|0-x\|=\left\|0-\mathrm{R}_{U} x\right\|=\left\|0-\mathrm{R}_{U \perp} x\right\|$,
which means that $(\forall x \in \mathcal{H}) C C_{S}(x)=0$.
Example 3.15. Assume that $\mathcal{H}=\mathbb{R}^{2}$ and $\mathcal{S}=\left\{T_{1}, T_{2}, T_{3}\right\}$, where for every $(x, y) \in$ $\mathbb{R}^{2}$,
$$
T_{1}(x, y)=(x, y) ; \quad T_{2}(x, y)=(-x, y) ; \quad T_{3}(x, y)=\left(x,-\frac{1}{4}(x-2)\right)
$$

Then
(i) $T_{1}(x, y), T_{2}(x, y), T_{3}(x, y)$ are affinely independent $\Leftrightarrow 2 x\left(-\frac{1}{4}(x-2)-y\right) \neq 0$;
(ii) $\left(\forall(x, y) \in \mathbb{R}^{2}\right) \quad C C_{\mathcal{S}}(x, y)=\left(0, \frac{1}{2}\left(y-\frac{1}{4}(x-2)\right)\right)$.

Consequently, $C C_{\mathcal{S}}$ is proper and continuous.
The following example shows that even if the operators in $\mathcal{S}$ are continuous, we generally have

$$
C C_{\mathcal{S}} \text { is proper } \nRightarrow C C_{\mathcal{S}} \text { is continuous. }
$$

Example 3.16. Assume that $\mathcal{H}=\mathbb{R}^{2}$ and $\mathcal{S}=\left\{T_{1}, T_{2}, T_{3}\right\}$, where for every $(x, y) \in$ $\mathbb{R}^{2}$,

$$
T_{1}(x, y)=(2,0) ; \quad T_{2}(x, y)=(-2,0) ; \quad T_{3}(x, y)=\left(x,-\frac{1}{4}(x-2)\right)
$$

Then
(i) $C C_{\mathcal{S}}$ is proper;
(ii) Let $(\forall k \geq 1)\left(x^{(k)}, y^{(k)}\right)=\left(2-\frac{1}{k}, 0\right)$. Then $C C_{\mathcal{S}}\left(x^{(k)}, y^{(k)}\right) \rightarrow(0,-8) \neq$ $(0,0)=C C_{\mathcal{S}}(2,0)$. Consequently, $C C_{\mathcal{S}}$ is not continuous at the point $(2,0)$.
Proof. (i): Let $(x, y) \in \mathbb{R}^{2}$. Now by Fact 2.7,
$T_{1}(x, y), T_{2}(x, y), T_{3}(x, y)$ are affinely independent $\Longleftrightarrow T_{2}(x, y)-T_{1}(x, y), T_{3}(x, y)-T_{1}(x, y)$ are linearly independent $\Longleftrightarrow(-4,0),\left(x-2,-\frac{1}{4}(x-2)\right)$ are linearly independent

$$
\begin{aligned}
& \Longleftrightarrow \operatorname{det}(A) \neq 0, \text { where } A=\left(\begin{array}{cc}
-4 & x-2 \\
0 & -\frac{1}{4}(x-2)
\end{array}\right) \\
& \Longleftrightarrow x-2 \neq 0 .
\end{aligned}
$$

Hence, by Corollary 2.15, when $x-2 \neq 0$, we have $C C_{\mathcal{S}}(x, y) \in \mathcal{H}$. Actually, when $x-2=0$, that is $x=2$, then for every $y \in \mathbb{R}$,

$$
T_{1}(2, y)=(2,0), T_{2}(2, y)=(-2,0), T_{3}(2, y)=\left(2,-\frac{1}{4}(2-2)\right)=(2,0)
$$

By Proposition 3.3, we know that $C C_{\mathcal{S}}(x, y)=(0,0) \in \mathcal{H}$. Hence, $C C_{\mathcal{S}}$ is proper.
(ii): Let $(\bar{x}, \bar{y})=(2,0)$, and $(\forall k \geq 1)\left(x^{(k)}, y^{(k)}\right)=\left(2-\frac{1}{k}, 0\right)$. By the analysis in (i) above, we know

$$
\begin{equation*}
C C_{\mathcal{S}}(\bar{x}, \bar{y})=(0,0) \tag{3.11}
\end{equation*}
$$

On the other hand, since

$$
\begin{aligned}
\mathcal{S}\left(x^{(k)}, y^{(k)}\right) & =\left\{T_{1}\left(x^{(k)}, y^{(k)}\right), T_{2}\left(x^{(k)}, y^{(k)}\right), T_{3}\left(x^{(k)}, y^{(k)}\right)\right\} \\
& =\left\{(2,0),(-2,0),\left(2-\frac{1}{k}, \frac{1}{4 k}\right)\right\},
\end{aligned}
$$

and since, by Definition 3.1, $C C_{\mathcal{S}}\left(x^{(k)}, y^{(k)}\right)=C C\left(\mathcal{S}\left(x^{(k)}, y^{(k)}\right)\right)$, we deduce that, by Fact 2.21 ,

$$
\begin{equation*}
C C_{\mathcal{S}}\left(x^{(k)}, y^{(k)}\right)=\left(0,-8+\frac{2}{k}+\frac{1}{8 k}\right) . \tag{3.12}
\end{equation*}
$$

Hence,

$$
\lim _{k \rightarrow \infty} C C_{\mathcal{S}}\left(x^{(k)}, y^{(k)}\right)=(0,-8) \neq(0,0) \stackrel{(3.11)}{=} C C_{\mathcal{S}}(2,0)
$$

and we are done.
3.4. The Demiclosedness Principle for circumcenter mappings. Let $T: \mathcal{H} \rightarrow$ $\mathcal{H}$ be nonexpansive. Then

$$
\left.\begin{array}{c}
x_{k} \rightharpoonup x  \tag{3.13}\\
x_{k}-T x_{k} \rightarrow 0
\end{array}\right\} \Rightarrow x \in \operatorname{Fix} T
$$

This well known implication (see [8, Theorem 3(a)]) is Browder's Demiclosedness Principle; it is a powerful tool in the study of nonexpansive mappings. (Technically speaking, (3.13) states that $\operatorname{Id}-T$ is demiclosed at 0 , but because a shift of a nonexpansive mapping is still nonexpansive, it is demiclosed everywhere.) For the sake of brevity, we shall simply say that
"the demiclosedness principle holds for $T$ " whenever (3.13) holds.
Clearly, the demiclosedness principle holds whenever $T$ is weak-to-strong continuous, e.g., when $T$ is continuous and $\mathcal{H}$ is finite-dimensional. The demiclosedness principle also holds for so-called subgradient projectors; see [5, Lemma 5.1] for details.

We now obtain a condition sufficient for the circumcenter mapping to satisfy the demiclosedness principle. Throughout, we assume $T_{1}, \ldots, T_{m}$ are mappings from $\mathcal{H}$ to $\mathcal{H}$.

Theorem 3.17. Suppose that the demiclosedness principle holds for each element in $\mathcal{S}=\left\{T_{1}, T_{2}, \ldots, T_{m}\right\}$. In addition, assume that $C C_{\mathcal{S}}$ is proper and that the implication

$$
\begin{equation*}
x_{k}-C C_{\mathcal{S}} x_{k} \rightarrow 0 \Rightarrow(\forall i \in\{1, \ldots, m\}) C C_{\mathcal{S}} x_{k}-T_{i} x_{k} \rightarrow 0 \tag{3.14}
\end{equation*}
$$

holds. Then the demiclosedness principle holds for $C C_{\mathcal{S}}$ and Fix $C C_{\mathcal{S}}=\bigcap_{i=1}^{m} \operatorname{Fix} T_{i}$.
Proof. Let $x_{k} \rightharpoonup x$ and

$$
\begin{equation*}
x_{k}-C C_{\mathcal{S}} x_{k} \rightarrow 0 \tag{3.15}
\end{equation*}
$$

By (3.15) and (3.14),

$$
(\forall i \in\{1, \ldots, m\}) \quad C C_{\mathcal{S}} x_{k}-T_{i} x_{k} \rightarrow 0
$$

Hence,

$$
(\forall i \in\{1, \ldots, m\}) \quad\left\|x_{k}-T_{i} x_{k}\right\| \leq\left\|x_{k}-C C_{\mathcal{S}} x_{k}\right\|+\left\|C C_{\mathcal{S}} x_{k}-T_{i} x_{k}\right\| \rightarrow 0
$$

Because the demiclosedness principle holds for each $T_{i}$, we deduce that $x \in \cap_{i=1}^{m} \operatorname{Fix} T_{i} \subseteq$ Fix $C C_{\mathcal{S}}$, where the last inclusion follows from Proposition 3.10(i). Therefore, $x-C C_{\mathcal{S}} x=0$, which shows that the demiclosedness principle holds for $C C_{\mathcal{S}}$. To verify the remaining assertion, let $\bar{x} \in \operatorname{Fix} C C_{\mathcal{S}}$. For every $k \in \mathbb{N}$, substitute $x_{k}$ by $\bar{x}$. Then using the assumption (3.14), we deduce that $(\forall i \in\{1, \ldots, m\})$ $C C_{\mathcal{S}} \bar{x}-T_{i} \bar{x}=0$. Combining with $(\forall i \in\{1, \ldots, m\})\left\|\bar{x}-T_{i} \bar{x}\right\| \leq\left\|\bar{x}-C C_{\mathcal{S}} \bar{x}\right\|+$ $\left\|C C_{\mathcal{S}} \bar{x}-T_{i} \bar{x}\right\|$, we obtain that $\bar{x} \in \cap_{i=1}^{m} \operatorname{Fix} T_{i}$. Hence, Fix $C C_{\mathcal{S}} \subseteq \cap_{i=1}^{m} \operatorname{Fix} T_{i}$. Therefore, the desired result follows from Proposition 3.10(i).

Corollary 3.18. Suppose that $T_{1}=\mathrm{Id}$ and that $C C_{\mathcal{S}}$ is proper. Then the implication

$$
\begin{equation*}
x_{k}-C C_{\mathcal{S}} x_{k} \rightarrow 0 \Rightarrow(\forall i \in\{1, \ldots, m\}) C C_{\mathcal{S}} x_{k}-T_{i} x_{k} \rightarrow 0 . \tag{3.16}
\end{equation*}
$$

holds.
Proof. Since $C C_{\mathcal{S}}$ is proper, by Remark 3.2, $\left\|C C_{\mathcal{S}} x_{k}-x_{k}\right\|=\left\|C C_{\mathcal{S}} x_{k}-T_{2} x_{k}\right\|=$ $\cdots=\left\|C C_{\mathcal{S}} x_{k}-T_{m} x_{k}\right\|$, which implies that (3.16) is true.
Proposition 3.19. Suppose that $T_{1}=\mathrm{Id}$, that for every $i \in\{2, \ldots, m\}$, the demiclosedness principle holds for $T_{i}$, that $\mathcal{S}=\left\{T_{1}, T_{2}, \ldots, T_{m}\right\}$, and that $C C_{\mathcal{S}}$ is proper. Then the demiclosedness principle holds for $C C_{\mathcal{S}}$ and $\operatorname{Fix} C C_{\mathcal{S}}=\cap_{i=1}^{m} \operatorname{Fix} T_{i}$.
Proof. Combine Theorem 3.17 with Corollary 3.18.
We are now ready for the main result of this section.
Theorem 3.20 (a demiclosedness principle for circumcenter mappings). Suppose that $T_{1}=\mathrm{Id}$, that each operator in $\mathcal{S}=\left\{T_{1}, T_{2}, \ldots, T_{m}\right\}$ is nonexpansive, and that $C C_{\mathcal{S}}$ is proper. Then the demiclosedness principle holds for $C C_{\mathcal{S}}$ and $\operatorname{Fix} C C_{\mathcal{S}}=$ $\cap_{i=1}^{m} \operatorname{Fix} T_{i}$.
Proof. Combine Browder's Demiclosedness Principle with Proposition 3.19.
We now present (omitting its easy proof) another consequence of Proposition 3.19.
Corollary 3.21. Suppose $\mathcal{H}$ is finite-dimensional, and $\mathcal{S}=\left\{T_{1}, \ldots, T_{m}\right\}$, where $T_{1}=\operatorname{Id}$ and $T_{j}$ is continuous for every $j \in\{2, \ldots, m\}$, and that $C C_{\mathcal{S}}$ is proper. Then the demiclosedness principle holds for $C C_{\mathcal{S}}$. In particular,

$$
\left.\begin{array}{c}
x_{k} \rightarrow \bar{x}  \tag{3.17}\\
x_{k}-C C_{\mathcal{S}} x_{k} \rightarrow 0
\end{array}\right\} \Rightarrow \bar{x} \in \cap_{j=1}^{m} \operatorname{Fix} T_{j}=\operatorname{Fix} C C_{\mathcal{S}}
$$

We now provide an example where the demiclosedness principle does not hold for $C C_{\mathcal{S}}$.
Example 3.22. Suppose that $\mathcal{H}=\mathbb{R}^{2}$. Set $L=\left\{(u, v) \in \mathcal{H} \left\lvert\, v=-\frac{1}{4} u+\frac{1}{2}\right.\right\}$. Assume that $\mathcal{S}=\left\{T_{1}, T_{2}, T_{3}\right\}$, where

$$
(\forall(u, v) \in \mathcal{H}) \quad T_{1}(u, v)=(-2,0), T_{2}(u, v)=(2,0) \text { and } T_{3}(u, v)=\mathrm{P}_{L}(u, v) .
$$

Set $\bar{x}=(0,-8)$ and $(\forall k \in \mathbb{N} \backslash\{0\}) x_{k}=\left(\frac{1}{k},-\frac{1}{4 k}-8\right)$. Then the following hold.
(i) $C C_{\mathcal{S}}$ is proper.
(ii) $\operatorname{Fix} C C_{\mathcal{S}}=\varnothing$.
(iii) $\lim _{k \rightarrow \infty} C C_{\mathcal{S}} x_{k}=\bar{x}=\lim _{k \rightarrow \infty} x_{k}$; consequently, $\lim _{k \rightarrow \infty}\left(x_{k}-C C_{\mathcal{S}} x_{k}\right)=0$. (See also Figure 1.)
(iv) $\bar{x} \notin$ Fix $C C_{\mathcal{S}}$; consequently, the demiclosedness principle does not hold for $C C_{\mathcal{S}}$.
Proof. (i): Let $x \in \mathcal{H}$. If $T_{3} x \in \mathbb{R} \cdot(1,0)$, then $T_{3} x=(2,0)$ and so $C C_{\mathcal{S}} x=(0,0)$. Now assume that $T_{3} x \notin \mathbb{R} \cdot(1,0)$. Then $T_{1} x, T_{2} x, T_{3} x$ are affinely independent. Hence, by Theorem 3.9, $C C_{\mathcal{S}} x \in \mathcal{H}$. Altogether, $C C_{\mathcal{S}}$ is proper.
(ii): Since $T_{1} x=(-2,0)$ and $T_{2} x=(2,0)$, by definition of circumcenter mapping,

$$
C C_{\mathcal{S}} x \in \mathbb{R} \cdot(0,1)
$$

which implies if $x \in \operatorname{Fix} C C_{\mathcal{S}}$, then $x \in \mathbb{R} \cdot(0,1)$. Since $T_{3}(0,-8)=\mathrm{P}_{L}(0,-8)=$ $(2,0)$, by Proposition 3.3, $C C_{\mathcal{S}}(0,-8)=(0,0) \neq(0,-8)$. Hence, $(0,-8) \notin \operatorname{Fix} C C_{\mathcal{S}}$. Let $x:=(0, v) \in(\mathbb{R} \cdot(0,1)) \backslash\{(0,-8)\}$. As seen in the proof of (i), the vectors $T_{1} x, T_{2} x, T_{3} x$ are affinely independent. Hence, by definition of circumcenter mapping, in this case $C C_{\mathcal{S}} x$ is the intersection of $\mathbb{R} \cdot(0,1)$ and the perpendicular bisector of the two points $T_{2} x, T_{3} x$. Denote by $C C_{\mathcal{S}} x:=(0, w)$. Some easy calculation yields that if $v>-8$, then $w>v$; if $v<-8$, then $w<v$, which means that $C C_{\mathcal{S}} x \neq x$. Altogether, $\operatorname{Fix} C C_{\mathcal{S}}=\varnothing$.
(iii): Let $k \in \mathbb{N} \backslash\{0\}$. Since $x_{k}=\left(\frac{1}{k},-\frac{1}{4 k}-8\right)$, by definition of $T_{3}$,

$$
T_{3} x_{k}=\left(2-\frac{1}{k}, \frac{1}{4 k}\right) .
$$

Hence

$$
C C_{\mathcal{S}} x_{k}=C C\left\{(-2,0),(2,0),\left(2-\frac{1}{k}, \frac{1}{4 k}\right)\right\} .
$$

By Example 3.16(ii), we obtain that

$$
\lim _{k \rightarrow \infty} C C_{\mathcal{S}} x_{k}=(0,-8)=\bar{x}=\lim _{k \rightarrow \infty} x_{k} .
$$

(iv): By (ii), $\bar{x} \notin$ Fix $C C_{\mathcal{S}}$. Therefore, the demiclosedness principle does not hold for $C C_{\mathcal{S}}$.


Figure 1. Example 3.22 illustrates that the demiclosedness principle may not hold for $C C_{\mathcal{S}}$.

Remark 3.23. Consider Example 3.22 where each $T_{i}$ is a projector and thus firmly nonexpansive but $\operatorname{Fix} C C_{\mathcal{S}}=\varnothing$. Is it possible to obtain an example where the demiclosedness principle does not hold but yet Fix $C C_{\mathcal{S}} \neq \varnothing$ ? We do not know the answer to this question.

## 4. Circumcenter mappings induced By Reflectors

Recall that $m \in \mathbb{N} \backslash\{0\}$ and that $U_{1}, \ldots, U_{m}$ are closed affine subspaces in the real Hilbert space $\mathcal{H}$ with $\cap_{i=1}^{m} U_{i} \neq \varnothing$. In the whole section, denote

$$
\Omega=\left\{\mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} \mid r \in \mathbb{N}, \text { and } i_{1}, i_{2}, \ldots, i_{r} \in\{1, \ldots, m\}\right\}
$$

By the empty product convention, $\prod_{j=1}^{0} \mathrm{R}_{U_{i_{j}}}=\mathrm{Id}$. Hence, $\Omega$ is a set that contains the identity operator, Id, and all (finite) compositions of the operators in $\left\{\mathrm{R}_{U_{1}}, \ldots, R_{U_{m}}\right\}$.
4.1. Proper circumcenter mappings induced by reflectors. Let us start with a useful lemma.

Lemma 4.1. Assume that $\mathrm{Id} \in \mathcal{S} \subseteq \Omega$. Let $x \in \mathcal{H}$. Then for every $\mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} \in$ $\mathcal{S}$,

$$
\left(\forall u \in \cap_{i=1}^{m} U_{i}\right) \quad\|x-u\|=\left\|\mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} x-u\right\|
$$

Proof. Let $u \in \cap_{i=1}^{m} U_{i}$. Because $U_{1}, \ldots, U_{m}$ are closed affine subspaces and $u \in$ $\cap_{i=1}^{m} U_{i} \subseteq \cap_{j=1}^{r} U_{i_{j}}$, by Proposition $2.10(\mathrm{v})$, we have

$$
\begin{aligned}
\|x-u\| & =\left\|\mathrm{R}_{U_{i_{1}}} x-u\right\| \\
\left\|\mathrm{R}_{U_{i_{1}}} x-u\right\| & =\left\|\mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} x-u\right\| \\
\vdots & \\
\left\|\mathrm{R}_{U_{i_{r-1}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} x-u\right\| & =\left\|\mathrm{R}_{U_{i_{r}}} \mathrm{R}_{U_{i_{r-1}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} x-u\right\|
\end{aligned}
$$

which yield

$$
\|x-u\|=\left\|\mathrm{R}_{U_{i_{r}}} \mathrm{R}_{U_{i_{r-1}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} x-u\right\|
$$

Proposition 4.2. Assume that $\operatorname{Id} \in \mathcal{S} \subseteq \Omega$. Let $x \in \mathcal{H}$. Then for every $u \in \cap_{i=1}^{m} U_{i}$,
(i) $\mathrm{P}_{\mathrm{aff}(\mathcal{S}(x))}(u) \in \operatorname{aff}(\mathcal{S}(x))$, and
(ii) for every $\mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} \in \mathcal{S}$,

$$
\left\|\mathrm{P}_{\mathrm{aff}(\mathcal{S}(x))}(u)-x\right\|=\left\|\mathrm{P}_{\mathrm{aff}(\mathcal{S}(x))}(u)-\mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} x\right\|
$$

Proof. (i): Let $u \in \cap_{i=1}^{m} U_{i}$. Because $\operatorname{aff}(\mathcal{S}(x))$ is the translate of a finite-dimensional linear subspace, $\operatorname{aff}(\mathcal{S}(x))$ is a closed affine subspace. Hence, we know $\mathrm{P}_{\text {aff }(\mathcal{S}(x))}(u)$ is well-defined. Clearly, $\mathrm{P}_{\mathrm{aff}(\mathcal{S}(x))}(u) \in \operatorname{aff}(\mathcal{S}(x))$, i.e., (i) is true.
(ii): Take $\mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}}$ in $\mathcal{S}$. Since Id, $\mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} \in \mathcal{S}$, we know $x, \mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} x \in \mathcal{S}(x) \subseteq \operatorname{aff}(\mathcal{S}(x))$. Denote $p=\mathrm{P}_{\text {aff }(\mathcal{S}(x))}(u)$. Substituting $C=\operatorname{aff}(\mathcal{S}(x)), x=u$ and $v=x$ in Proposition 2.10(iii), we deduce

$$
\begin{equation*}
\|u-p\|^{2}+\|x-p\|^{2}=\|u-x\|^{2} \tag{4.1}
\end{equation*}
$$

Similarly, substitute $C=\operatorname{aff}(\mathcal{S}(x)), x=u$ and $v=\mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} x$ in Proposition 2.10(iii) to obtain

$$
\begin{equation*}
\|u-p\|^{2}+\left\|\mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} x-p\right\|^{2}=\left\|u-\mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} x\right\|^{2} \tag{4.2}
\end{equation*}
$$

On the other hand, by Lemma 4.1, we know

$$
\begin{equation*}
\|x-u\|=\left\|\mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} x-u\right\| \tag{4.3}
\end{equation*}
$$

Combining (4.3) with (4.1) and (4.2), we yield

$$
\|p-x\|=\left\|p-\mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} x\right\|
$$

Since $\mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} \in \mathcal{S}$ is arbitrary, thus (ii) holds.
Combining Proposition 3.6 with Proposition 4.2 , we deduce the theorem below which is one of the main results in this paper.

Theorem 4.3. Assume that $\operatorname{Id} \in \mathcal{S} \subseteq \Omega$. Then the following hold:
(i) The circumcenter mapping $C C_{\mathcal{S}}: \mathcal{H} \rightarrow \mathcal{H}$ induced by $\mathcal{S}$ is proper, i.e., for every $x \in \mathcal{H}, C C_{\mathcal{S}} x$ is the unique point satisfying the two conditions below:
(a) $C C_{\mathcal{S}} x \in \operatorname{aff}(\mathcal{S}(x))$, and
(b) $\left(\forall \mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{1}}} \in \mathcal{S}\right)\left\|C C_{\mathcal{S}} x-x\right\|=\left\|C C_{\mathcal{S}} x-\mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{1}}} x\right\|$.
(ii) $(\forall x \in \mathcal{H})\left(\forall u \in \cap_{i=1}^{m} U_{i}\right) C C_{\mathcal{S}} x=\mathrm{P}_{\mathrm{aff}(\mathcal{S}(x))}(u)$.
(iii) $(\forall x \in \mathcal{H}) C C_{\mathcal{S}} x=\mathrm{P}_{\mathrm{aff}(\mathcal{S}(x))}\left(\mathrm{P}_{\cap_{i=1}^{m} U_{i}} x\right)$.

Proof. (i) and (ii): The required results follow from Propositions 3.6 and 4.2.
(iii): Since $\mathrm{P}_{\cap_{i=1}^{m} U_{i}} x \in \cap_{i=1}^{m} U_{i}$, the desired result comes from (ii).

We now list several proper circumcenter mappings induced by reflectors; the properness of some of these mappings is derived from Theorem 4.3.

Example 4.4. Assume that $\mathcal{S}=\left\{\operatorname{Id}, \mathrm{R}_{U_{1}}, \ldots, \mathrm{R}_{U_{m}}\right\}$. By Theorem 4.3(i), $C C_{\mathcal{S}}$ is proper.

Example 4.5. Assume that

$$
\mathcal{S}=\left\{\operatorname{Id}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{3}} \mathrm{R}_{U_{2}}, \ldots, \mathrm{R}_{U_{m}} \mathrm{R}_{U_{m-1}}, \mathrm{R}_{U_{1}} \mathrm{R}_{U_{m}}\right\}
$$

By Theorem 4.3(i), $C C_{\mathcal{S}}$ is proper.
Example 4.6 (Behling et al. [6]). Assume that $m=2$ and that $\mathcal{S}=\left\{\operatorname{Id}, \mathrm{R}_{U_{1}}\right.$, $\left.\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}\right\}$. Then, by Theorem $4.3(\mathrm{i}), C C_{\mathcal{S}}$ is proper.

Example 4.7 (Behling et al. [7]). Assume that

$$
\mathcal{S}=\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}, \ldots, \mathrm{R}_{U_{m}} \cdots \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}\right\}
$$

Then, by Theorem $4.3(\mathrm{i}), C C_{\mathcal{S}}$ is proper.
Remark 4.8. In fact, the C-DRM operator $C_{T}$ defined in [6, Section 2] is the $C C_{\mathcal{S}}$ operator of Example 4.6 while the CRM operator $C$ defined in [7, page 159] is the operator $C C_{\mathcal{S}}$ from Example 4.7.

Example 4.9. Assume that $m=2$ and that $\mathcal{S}=\left\{\mathrm{Id}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}\right\}$. By Proposition 3.3,

$$
C C_{\mathcal{S}}=\frac{\mathrm{Id}+\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}}{2}
$$

which is the well-known Douglas-Rachford splitting operator. Clearly, $C C_{\mathcal{S}}$ is proper.
Example 4.10. Assume that $m=2$ and that $\mathcal{S}=\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}}\right\}$. Then $C C_{\mathcal{S}}$ is proper. Moreover,

$$
\left(\forall x \in U_{1}\right) \quad C C_{\mathcal{S}} x=\mathrm{P}_{U_{2}} x \quad \text { and } \quad\left(\forall x \in U_{2}\right) \quad C C_{\mathcal{S}} x=\mathrm{P}_{U_{1}} x
$$

Proof. The first assertion follows from Example 4.4. As for the remaining ones, note that

$$
\begin{equation*}
\left(\forall x \in U_{1}\right) \mathcal{S}(x)=\left\{x, \mathrm{R}_{U_{2}} x\right\} \text { and }\left(\forall x \in U_{2}\right) \mathcal{S}(x)=\left\{x, \mathrm{R}_{U_{1}} x\right\} \tag{4.4}
\end{equation*}
$$

Combining (4.4) with Proposition 3.3, we obtain that

$$
\begin{aligned}
& \left(\forall x \in U_{1}\right) \quad C C_{\mathcal{S}} x=\frac{x+\mathrm{R}_{U_{2}} x}{2}=\mathrm{P}_{U_{2}} x \quad \text { and } \\
& \left(\forall x \in U_{2}\right) \quad C C_{\mathcal{S}} x=\frac{x+\mathrm{R}_{U_{1}} x}{2}=\mathrm{P}_{U_{1}} x
\end{aligned}
$$

The proof is complete.
Example 4.11. Assume that $m=2$ and that $\mathcal{S}=\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}\right\}$. Let $x \in$ $\mathcal{H}$ and set $l=\operatorname{card}\left\{x, \mathrm{R}_{U_{1}} x, \mathrm{R}_{U_{2}} x, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} x\right\}$. Then exactly one of the following cases occurs.
(i) $l=1$ and $C C_{\mathcal{S}} x=x$.
(ii) $l=2$, say $S(x)=\left\{x_{1}, x_{2}\right\}$, where $x_{1}$ and $x_{2}$ are two distinct elements in $S(x)$, and $C C_{\mathcal{S}} x=\frac{x_{1}+x_{2}}{2}$.
(iii) $l=3$, say $S(x)=\left\{x_{1}, x_{2}, x_{3}\right\}$, where $x_{1}, x_{2}, x_{3}$ are pairwise distinct elements in $S(x)$, and $C C_{\mathcal{S}} x=\frac{N}{D}$ where $N:=\left\|x_{2}-x_{3}\right\|^{2}\left\langle x_{1}-x_{3}, x_{1}-\right.$ $\left.x_{2}\right\rangle x_{1}+\left\|x_{1}-x_{3}\right\|^{2}\left\langle x_{2}-x_{3}, x_{2}-x_{1}\right\rangle x_{2}+\left\|x_{1}-x_{2}\right\|^{2}\left\langle x_{3}-x_{1}, x_{3}-x_{2}\right\rangle x_{3}$ and $D:=2\left(\left\|x_{2}-x_{1}\right\|^{2}\left\|x_{3}-x_{1}\right\|^{2}-\left\langle x_{2}-x_{1}, x_{3}-x_{1}\right\rangle^{2}\right)$.
(iv) $l=4$ and

$$
\begin{gathered}
C C_{\mathcal{S}} x=x_{1}+\frac{1}{2} \mathbf{u} G\left(x_{2}-x_{1}, \ldots, x_{t_{x}-1}-x_{1}, x_{t_{x}}-x_{1}\right)^{-1} \mathbf{v} \\
\text { where }\left\{x_{1}, \ldots, x_{t_{x}}\right\} \quad \subseteq \\
\operatorname{dim}(\operatorname{aff} S(x)) \text { and } x_{1}, x_{2}, \ldots, x_{t_{x}} \text { are affinely independent, and } \mathbf{u}:=\left(x_{2}-\right. \\
\left.x_{1}, \ldots, x_{t_{x}}-x_{1}\right), \mathbf{v}:=\left(\left\|x_{2}-x_{1}\right\|^{2}, \ldots,\left\|x_{t_{x}}-x_{1}\right\|^{2}\right)^{T} \text { and } \\
G\left(x_{2}-x_{1}, \ldots, x_{t_{x}-1}-x_{1}, x_{t_{x}}-x_{1}\right) \\
=\left(\begin{array}{cccc}
\left\|x_{2}-x_{1}\right\|^{2} & \left\langle x_{2}-x_{1}, x_{3}-x_{1}\right\rangle & \ldots & \left\langle x_{2}-x_{1}, x_{t_{x}}-x_{1}\right\rangle \\
\vdots & \vdots & & \vdots \\
\left\langle x_{t_{x}-1}-x_{1}, x_{2}-x_{1}\right\rangle & \left\langle x_{t_{x}-1}-x_{1}, x_{3}-x_{1}\right\rangle & \ldots & \left\langle x_{t_{x}-1}-x_{1}, x_{t_{x}}-x_{1}\right\rangle \\
\left\langle x_{t_{x}}-x_{1}, x_{2}-x_{1}\right\rangle & \left\langle x_{t_{x}}-x_{1}, x_{3}-x_{1}\right\rangle & \ldots & \left\|x_{t_{x}}-x_{1}\right\|^{2}
\end{array}\right) .
\end{gathered}
$$

Proof. By Theorem 4.3(i), $C C_{\mathcal{S}}$ is proper. The rest follows from Facts 2.12 and 2.13.

We now turn to the properness of $C C_{\widetilde{\mathcal{S}}}$ when $\operatorname{Id} \in \widetilde{\mathcal{S}} \subseteq$ aff $\Omega$.
Proposition 4.12. Let $\alpha \in \mathbb{R}$. Assume that

$$
\begin{equation*}
\widetilde{\mathcal{S}}=\left\{\operatorname{Id},(1-\alpha) \operatorname{Id}+\alpha \mathrm{R}_{U_{1}}, \ldots,(1-\alpha) \operatorname{Id}+\alpha \mathrm{R}_{U_{m}}\right\} \tag{4.5}
\end{equation*}
$$

and that

$$
\begin{equation*}
\mathcal{S}=\left\{\operatorname{Id}, \mathrm{R}_{U_{1}}, \ldots, \mathrm{R}_{U_{m}}\right\} . \tag{4.6}
\end{equation*}
$$

Then $C C_{\tilde{\mathcal{S}}}$ is proper. Moreover,

$$
\begin{equation*}
(\forall x \in \mathcal{H}) \quad C C_{\widetilde{\mathcal{S}}} x=\alpha C C_{\mathcal{S}} x+(1-\alpha) x \in \mathcal{H} \tag{4.7}
\end{equation*}
$$

Proof. If $\alpha=0$, then $\widetilde{\mathcal{S}}=\{\mathrm{Id}\}$, by Definition 3.1,

$$
(\forall x \in \mathcal{H}) \quad C C_{\widetilde{\mathcal{S}}} x=x=0 C C_{\mathcal{S}} x+(1-0) x \in \mathcal{H} .
$$

Now assume $\alpha \neq 0$. Let $x \in \mathcal{H}$. For every $i \in\{1, \ldots, m\}$, thus

$$
\begin{aligned}
& C C_{\widetilde{\mathcal{S}}} x \\
= & C C(\widetilde{\mathcal{S}}(x)) \quad \text { (by Definition 3.1) } \\
= & C C\left(\left\{x,(1-\alpha) x+\alpha \mathrm{R}_{U_{1}} x, \ldots,(1-\alpha) x+\alpha \mathrm{R}_{U_{m}} x\right\}\right) \quad(\text { by (4.5)) } \\
= & C C\left(\left\{0, \alpha\left(\mathrm{R}_{U_{1}} x-x\right), \ldots, \alpha\left(\mathrm{R}_{U_{m}} x-x\right)\right\}+x\right) \\
= & C C\left(\left\{0, \alpha\left(\mathrm{R}_{U_{1}} x-x\right), \ldots, \alpha\left(\mathrm{R}_{U_{m}} x-x\right)\right\}\right)+x \quad \text { (by Fact 2.18) } \\
= & \left.\alpha C C\left(\left\{0, \mathrm{R}_{U_{1}} x-x, \ldots, \mathrm{R}_{U_{m}} x-x\right\}\right)+x \quad \quad \text { by Fact } 2.17 \text { and } \alpha \neq 0\right) \\
= & \alpha C C\left(\left\{x, \mathrm{R}_{U_{1}} x, \ldots, \mathrm{R}_{U_{m}} x\right\}-x\right)+x \\
= & \alpha C C\left(\left\{x, \mathrm{R}_{U_{1}} x, \ldots, \mathrm{R}_{U_{m}} x\right\}\right)-\alpha x+x \quad \text { (by Fact 2.18) } \\
= & \alpha C C(\mathcal{S}(x))+(1-\alpha) x \quad \text { (by (4.6)) } \\
= & \alpha C C_{\mathcal{S}} x+(1-\alpha) x \in \mathcal{H} . \quad \text { (by Definition 3.1 and Theorem 4.3(i)) }
\end{aligned}
$$

The proof is complete.
Proposition 4.13. Assume that $\mathcal{S}=\left\{\mathrm{Id}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}\right\}$, set $T=$ $\frac{\mathrm{Id}+\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}}{2}$, which is the Douglas-Rachford splitting operator, and set $\widetilde{\mathcal{S}}=\left\{\mathrm{Id}, T, T^{2}\right\}$. Then the following hold:
(i) $\operatorname{aff}\left\{\operatorname{Id}, T, T^{2}\right\}=\operatorname{aff} \mathcal{S}$.
(ii) $C C_{\widetilde{\mathcal{S}}}$ is proper.

Proof. (i): By Fact 2.5, aff $\left\{\operatorname{Id}, T, T^{2}\right\}=\operatorname{aff}(\mathcal{S})$ if and only if $\operatorname{Id}+\operatorname{span}\left\{T-\mathrm{Id}, T^{2}-\right.$ $\mathrm{Id}\}=\mathrm{Id}+\operatorname{span}\left\{\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}-\mathrm{Id}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}-\mathrm{Id}\right\}$. On the other hand,

$$
\begin{equation*}
T-\mathrm{Id}=\frac{\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}+\mathrm{Id}}{2}-\mathrm{Id}=\frac{\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}-\mathrm{Id}}{2}, \tag{4.8}
\end{equation*}
$$

and

$$
\begin{align*}
T^{2}-\mathrm{Id} & =T^{2}-T+T-\mathrm{Id}=(T-\mathrm{Id}) T+(T-\mathrm{Id})  \tag{4.9}\\
& =\frac{\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}-\mathrm{Id}}{2}\left(\frac{\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}+\mathrm{Id}}{2}\right)+\frac{\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}-\mathrm{Id}}{2} \\
& =\frac{1}{4}\left(\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}-\mathrm{Id}\right)+\frac{1}{2}\left(\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}-\mathrm{Id}\right),
\end{align*}
$$

which result in

$$
\left(T-\mathrm{Id}, \quad T^{2}-\mathrm{Id}\right)=\left(\begin{array}{l}
\left.\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}-\mathrm{Id}, \quad \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}-\mathrm{Id}\right)\left(\begin{array}{cc}
\frac{1}{2} & \frac{1}{2} \\
0 & \frac{1}{4}
\end{array}\right) . . . . . \tag{4.10}
\end{array}\right.
$$

Set $A=\left(\begin{array}{cc}\frac{1}{2} & \frac{1}{2} \\ 0 & \frac{1}{4}\end{array}\right)$. Since $\operatorname{det}(A)=\frac{1}{8} \neq 0,(4.10)$ yields

$$
\operatorname{span}\left\{T-\mathrm{Id}, T^{2}-\mathrm{Id}\right\}=\operatorname{span}\left\{\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}-\mathrm{Id}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}-\mathrm{Id}\right\}
$$

Altogether, the proof of (i) is complete.
(ii): If $x, T x, T^{2} x$ are affinely independent, by Fact 2.13 , then $C C_{\widetilde{\mathcal{S}}} x \in \mathcal{H}$. Suppose $x, T x, T^{2} x$ are affinely dependent. By (4.10) above and $\operatorname{det}(A) \neq 0$, in this case, $x, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} x, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} x$ are affinely dependent. Applying Theorem 4.3(i), we know $C C_{\mathcal{S}} x \in \mathcal{H}$. Hence, Fact 2.14 yields that

$$
\begin{equation*}
\operatorname{card}\left(\left\{x, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}, x \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} x\right\}\right)=\operatorname{card}(\mathcal{S}(x)) \leq 2 \tag{4.11}
\end{equation*}
$$

If $T x-x=0$, by Proposition 3.3, $C C_{\widetilde{\mathcal{S}}} x=\frac{x+T^{2} x}{2}$. Now suppose $T x-x \neq 0$. By (4.8), $\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} x \neq x$. Therefore, by (4.11) and $\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} x \neq x$, either $\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} x=$ $\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} x$ or $\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} x=x$. Suppose $\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} x=\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} x$. Multiply both sides by $\mathrm{R}_{U_{1}} \mathrm{R}_{U_{2}}$, by Lemma 2.8(ii), to deduce $\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} x=x$, which contradicts with $\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} x \neq x$. Suppose $\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}} x=x$, by (4.8) and (4.9), which implies, $T x=T^{2} x$. Then by Proposition 3.3, we obtain $C C_{\widetilde{\mathcal{S}}} x=\frac{x+T x}{2} \in \mathcal{H}$.

In conclusion, $(\forall x \in \mathcal{H}) C C_{\widetilde{\mathcal{S}}} x \in \mathcal{H}$, which means (ii) holds.

### 4.2. Improper circumcenter mappings induced by reflectors.

Propositions 4.12 and 4.13 naturally prompt the following question: Is $C C_{\widetilde{\mathcal{S}}}$ proper for every $\widetilde{S}$ with Id $\in \widetilde{S} \subseteq$ aff $\Omega$ ? The following examples provide negative answers.

Example 4.14. Assume that $m=2$, that $U:=U_{1}=U_{2} \varsubsetneqq \mathcal{H}$, and that $\left\{\alpha_{1}, \alpha_{2}\right\} \subseteq$ $\mathbb{R}$. Assume further that $\widetilde{\mathcal{S}}=\left\{\operatorname{Id},\left(1-\alpha_{1}\right) \operatorname{Id}+\alpha_{1} \mathrm{R}_{U},\left(1-\alpha_{2}\right) \operatorname{Id}+\alpha_{2} \mathrm{R}_{U}\right\}$. Then $C C_{\widetilde{\mathcal{S}}}$ is improper if and only if $\alpha_{1} \neq 0, \alpha_{2} \neq 0$ and $\alpha_{2} \neq \alpha_{1}$.

Proof. By Proposition 3.3, when $\alpha_{1}=0$ or $\alpha_{2}=0$, then $C C_{\widetilde{\mathcal{S}}}$ is proper.
For every $x \in \mathcal{H}$, if $\alpha_{1} \neq 0$,

$$
\begin{aligned}
\left(\left(1-\alpha_{2}\right) x+\alpha_{2} \mathrm{R}_{U} x\right)-x & =\alpha_{2}\left(\mathrm{R}_{U} x-x\right) \\
& =\frac{\alpha_{2}}{\alpha_{1}} \alpha_{1}\left(\mathrm{R}_{U} x-x\right) \\
& =\frac{\alpha_{2}}{\alpha_{1}}\left(\left(\left(1-\alpha_{1}\right) x+\alpha_{1} \mathrm{R}_{U} x\right)-x\right)
\end{aligned}
$$

which implies that, by Fact 2.7,

$$
x,\left(1-\alpha_{1}\right) x+\alpha_{1} \mathrm{R}_{U} x,\left(1-\alpha_{2}\right) x+\alpha_{2} \mathrm{R}_{U} x \text { are affinely dependent. }
$$

On the other hand, if $x \in \mathcal{H} \backslash U$, then since $\alpha_{1} \neq 0, \alpha_{2} \neq 0, \alpha_{1} \neq \alpha_{2}$ and $\operatorname{Fix}_{U}=U$, we obtain that

$$
\begin{aligned}
& \operatorname{card}\left\{x,\left(1-\alpha_{1}\right) x+\alpha_{1} \mathrm{R}_{U} x,\left(1-\alpha_{2}\right) x+\alpha_{2} \mathrm{R}_{U} x\right\}=3 \\
\Longleftrightarrow & \alpha_{1} \neq 0, \alpha_{2} \neq 0 \text { and } \alpha_{2} \neq \alpha_{1}
\end{aligned}
$$

Therefore, we deduce the required result.
Example 4.15. Assume that $m=2$, that $U:=U_{1}=U_{2} \varsubsetneqq \mathcal{H}$, and that $\left\{\alpha_{1}, \alpha_{2}\right\} \subseteq$ $\mathbb{R}$. Assume further that $\widetilde{\mathcal{S}}=\left\{\operatorname{Id},\left(1-\alpha_{1}\right) \operatorname{Id}+\alpha_{1} \mathrm{R}_{U},\left(\left(1-\alpha_{2}\right) \operatorname{Id}+\alpha_{2} \mathrm{R}_{U}\right) \circ((1-\right.$ $\left.\left.\left.\alpha_{1}\right) \operatorname{Id}+\alpha_{1} \mathrm{R}_{U}\right)\right\}$. Then $C C_{\widetilde{\mathcal{S}}}$ is improper if and only if $\alpha_{1} \neq 0, \alpha_{1} \neq \frac{1}{2}, \alpha_{2} \neq 0$ and $\alpha_{2} \neq \frac{\alpha_{1}}{2 \alpha_{1}-1}$.
Proof. By Proposition 3.3, when $\alpha_{1}=0$ or $\alpha_{2}=0$, then $C C_{\widetilde{\mathcal{S}}}$ is proper.
Note that

$$
\begin{aligned}
& \left(\left(1-\alpha_{2}\right) \operatorname{Id}+\alpha_{2} \mathrm{R}_{U}\right) \circ\left(\left(1-\alpha_{1}\right) \operatorname{Id}+\alpha_{1} \mathrm{R}_{U}\right) \\
= & \left(1-\alpha_{1}-\alpha_{2}+2 \alpha_{2} \alpha_{1}\right) \operatorname{Id}+\left(\alpha_{1}+\alpha_{2}-2 \alpha_{2} \alpha_{1}\right) \mathrm{R}_{U} \quad \text { (by Lemma 2.8(ii)). }
\end{aligned}
$$

Hence, for every $x \in \mathcal{H}$, if $\alpha_{1} \neq 0$,

$$
\begin{aligned}
& \left(\left(\left(1-\alpha_{2}\right) \operatorname{Id}+\alpha_{2} \mathrm{R}_{U}\right) \circ\left(\left(1-\alpha_{1}\right) \mathrm{Id}+\alpha_{1} \mathrm{R}_{U}\right) x\right)-x \\
& =\left(\alpha_{1}+\alpha_{2}-2 \alpha_{2} \alpha_{1}\right)\left(\mathrm{R}_{U} x-x\right) \\
& =\frac{\alpha_{1}+\alpha_{2}-2 \alpha_{2} \alpha_{1}}{\alpha_{1}} \alpha_{1}\left(\mathrm{R}_{U} x-x\right) \\
& =\frac{\alpha_{1}+\alpha_{2}-2 \alpha_{2} \alpha_{1}}{\alpha_{1}}\left(\left(\left(1-\alpha_{1}\right) x+\alpha_{1} \mathrm{R}_{U} x\right)-x\right)
\end{aligned}
$$

which implies, by Fact 2.7 , that $x,\left(1-\alpha_{1}\right) x+\alpha_{1} \mathrm{R}_{U} x,\left(\left(1-\alpha_{2}\right) \operatorname{Id}+\alpha_{2} \mathrm{R}_{U}\right) \circ((1-$ $\left.\left.\alpha_{1}\right) \operatorname{Id}+\alpha_{1} \mathrm{R}_{U}\right) x$ are affinely dependent. On the other hand, assume now $x \in \mathcal{H} \backslash U$. Then card $\left\{x,\left(1-\alpha_{1}\right) x+\alpha_{1} \mathrm{R}_{U} x,\left(\left(1-\alpha_{2}\right) \operatorname{Id}+\alpha_{2} \mathrm{R}_{U}\right) \circ\left(\left(1-\alpha_{1}\right) \operatorname{Id}+\alpha_{1} \mathrm{R}_{U}\right) x\right\}=3$ if and only if $\alpha_{1} \neq 0, \alpha_{1} \neq \frac{1}{2}, \alpha_{2} \neq 0$ and $\alpha_{2} \neq \frac{\alpha_{1}}{2 \alpha_{1}-1}$. Therefore, combining Corollary 3.4 with the results obtained above, we infer the desired result.

The following example is a special case of Example 4.15.
Example 4.16. Assume that $m=2$, that $U_{1} \varsubsetneqq U_{2}=\mathcal{H}$, and that $\left\{\alpha_{1}, \alpha_{2}\right\} \subseteq \mathbb{R}$. Assume further that $\widetilde{\mathcal{S}}=\left\{\operatorname{Id},\left(1-\alpha_{1}\right) \operatorname{Id}+\alpha_{1} \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}},\left(\left(1-\alpha_{2}\right) \mathrm{Id}+\alpha_{2} \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}\right) \circ\right.$ $\left.\left(\left(1-\alpha_{1}\right) \operatorname{Id}+\alpha_{1} \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}\right)\right\}$. Then $C C_{\widetilde{\mathcal{S}}}$ is improper if and only if $\alpha_{1} \neq 0, \alpha_{1} \neq$ $\frac{1}{2}, \alpha_{2} \neq 0$ and $\alpha_{2} \neq \frac{\alpha_{1}}{2 \alpha_{1}-1}$.
Proof. Since $\mathrm{R}_{U_{2}}=\mathrm{R}_{\mathcal{H}}=\mathrm{Id}$, we deduce that

$$
\widetilde{\mathcal{S}}=\left\{\operatorname{Id},\left(1-\alpha_{1}\right) \operatorname{Id}+\alpha_{1} \mathrm{R}_{U_{1}},\left(\left(1-\alpha_{2}\right) \operatorname{Id}+\alpha_{2} \mathrm{R}_{U_{1}}\right)\left(\left(1-\alpha_{1}\right) \operatorname{Id}+\alpha_{1} \mathrm{R}_{U_{1}}\right)\right\}
$$

The desired result follows directly from Example 4.15.

Notice that in Proposition 4.13 we showed that for $\widetilde{\mathcal{S}}=\left\{\operatorname{Id}, T, T^{2}\right\}=$ $\left\{\mathrm{Id}, \frac{\mathrm{Id}+\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}}{2}, \frac{\mathrm{Id}+\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}}{2} \circ \frac{\mathrm{Id}+\mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}}{2}\right\}, C C_{\widetilde{\mathcal{S}}}$ is proper. The example above says that this result is not a conincidence.

### 4.3. Particular circumcenter mappings in Euclidean spaces.

### 4.3.1. Application to best approximation. Suppose that

$$
\mathcal{S}_{1}=\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}}\right\} \quad \text { and } \quad \mathcal{S}_{2}=\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}\right\} .
$$

By Example 4.4 and Example 4.6, we know $C C_{S_{1}}$ and $C C_{S_{2}}$ are proper. Hence, for every $x \in \mathcal{H}$, we are able to generate iterations $\left(C C_{S_{1}}^{k} x\right)_{k \in \mathbb{N}}$ and $\left(C C_{S_{2}}^{k} x\right)_{k \in \mathbb{N}}$.

In the following two examples, we choose two linear subspaces, $U_{1}$ and $U_{2}$, in $\mathbb{R}^{3}$ and one point $x_{0} \in \mathbb{R}^{3}$. Then we count the iteration numbers needed for the four algorithms: the shadow sequence of the Douglas-Rachford method (DRM) (see, [1] for details), the sequence generated by the method of alternating projections (MAP), and the sequence generated by iterating $C C_{S_{1}}$ and $C C_{S_{2}}$ to find the best approximation point $\bar{x}=\mathrm{P}_{U_{1} \cap U_{2}} x_{0}$.

Example 4.17. Assume that $\mathcal{H}=\mathbb{R}^{3}$, that $U_{1}$ is the line passing through the points $(0,0,0)$ and $(1,0,0)$, and that $U_{2}$ is the plane $\{(x, y, z) \mid x+y+z=0\}$. Let $x_{0}=(0.5,0,0)$. As Table 1 shows, both of the $C C_{S_{1}}$ and $C C_{S_{2}}$ are faster than DRM and MAP. (The results were obtained using GeoGebra.)

| Algorithm | Iterations needed to find $\mathrm{P}_{U_{1} \cap U_{2}} x_{0}$ |
| :--- | :--- |
| Shadow DRM | 12 |
| MAP | 12 |
| Iterating $C C_{S_{1}}$ | 1 |
| Iterating $C C_{S_{2}}$ | 1 |

Table 1. Iterations required for each algorithm. See Example 4.17 for details.


Figure 2. Example 4.17 compares iterations for a line and a plane.

Example 4.18. Assume that $\mathcal{H}=\mathbb{R}^{3}$, that $U_{1}=\{(x, y, z) \mid x+y+z=0\}$, and that $U_{2}:=\{(x, y, z) \mid-x+2 y+2 z=0\}$. Set $x_{0}=(-1,0.5,0.5)$. As Table 2 illustrates, $C C_{S_{2}}$ is faster than the other methods, and $C C_{S_{1}}$ performs no worse than DRM or MAP. (The results were obtained using GeoGebra.)

| Algorithm | Iterations needed to find $\mathrm{P}_{U_{1} \cap U_{2}} x_{0}$ |
| :--- | :--- |
| Shadow DRM | 5 |
| MAP | 6 |
| Iterating $C C_{S_{1}}$ | 5 |
| Iterating $C C_{S_{2}}$ | 2 |

TAble 2. Iterations required for each algorithm. See Example 4.18 for details.


Figure 3. Example 4.18 compares iterations for two planes.
4.3.2. Counterexamples. The following two examples show that the circumcenter mapping induced by reflectors is in general neither linear nor continuous.

Example 4.19 (Discontinuity). Suppose that $\mathcal{H}=\mathbb{R}^{2}$, set $U_{1}=\mathbb{R} \cdot(1,0)$, and set $U_{2}:=\mathbb{R} \cdot(1,1)$. Suppose that $\mathcal{S}=\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}}\right\}$ or that $\mathcal{S}=\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}\right\}$. Let $\bar{x}=(1,0)$ and let $(\forall k \in \mathbb{N}) x_{k}=\left(1, \frac{1}{k+1}\right)$. As Figure 4 illustrates, $C C_{\mathcal{S}} \bar{x}=\left(\frac{1}{2}, \frac{1}{2}\right)$ and $(\forall k \in \mathbb{N}) C C_{\mathcal{S}} x_{k}=(0,0)$. Hence,

$$
\lim _{k \rightarrow \infty} C C_{\mathcal{S}} x_{k}=(0,0) \neq\left(\frac{1}{2}, \frac{1}{2}\right)=C C_{\mathcal{S}} \bar{x}
$$

which implies that $C C_{\mathcal{S}}$ is not continuous at $\bar{x}$. By Corollary 3.21, the demiclosedness principle holds for $C C_{\mathcal{S}}$.


Figure 4. Example 4.19 provides a discontinuous $C C_{\mathcal{S}}$ in $\mathbb{R}^{2}$.
Example 4.20 (Nonlinearity). Suppose that $\mathcal{H}=\mathbb{R}^{2}$, set $U_{1}=\mathbb{R} \cdot(1,0)$ and set $U_{2}=\mathbb{R} \cdot(1,1)$. Suppose that $\mathcal{S}=\left\{\operatorname{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}}\right\}$ or that $\mathcal{S}=\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}\right\}$. Let $x=(1,0)$ and $y=(1,-1)$. As Figure 5 illustrates,

$$
C C_{\mathcal{S}} x+C C_{\mathcal{S}} y=\left(\frac{1}{2}, \frac{1}{2}\right)+(0,0) \neq(0,0)=C C_{\mathcal{S}}(x+y)
$$

which shows that $C C_{\mathcal{S}}$ is not linear. By Corollary 3.21 , the demiclosedness principle holds for $C C_{\mathcal{S}}$.


Figure 5. Example 4.20 presents a nonlinear $C C_{\mathcal{S}}$ in $\mathbb{R}^{2}$
5. Circumcenter mappings induced By projectors

In this section, we uphold the notations that

$$
\Omega=\left\{\mathrm{R}_{U_{i_{r}}} \cdots \mathrm{R}_{U_{i_{2}}} \mathrm{R}_{U_{i_{1}}} \mid r \in \mathbb{N}, \text { and } i_{1}, \ldots, i_{r} \in\{1, \ldots, m\}\right\}
$$

and $\operatorname{Id} \in \mathcal{S} \subseteq \Omega$. In addition, set

$$
\Theta=\left\{\mathrm{P}_{U_{i_{r}}} \cdots \mathrm{P}_{U_{i_{2}}} \mathrm{P}_{U_{i_{1}}} \mid r \in \mathbb{N}, \text { and } i_{1}, \ldots, i_{r} \in\{1, \ldots, m\}\right\}
$$

By the empty product convention, $\prod_{j=1}^{0} \mathrm{P}_{U_{i_{j}}}=\mathrm{Id}$. Hence Id $\in \Theta$. Specifically, we assume that

$$
\operatorname{Id} \in \widehat{\mathcal{S}} \subseteq \operatorname{aff} \Theta
$$

5.1. Proper circumcenter mappings induced by projectors. First, we present some cases when $C C_{\widehat{\mathcal{S}}}$ is proper.
Proposition 5.1. Let $\alpha \in \mathbb{R}$. Assume that

$$
\widehat{\mathcal{S}}=\left\{\operatorname{Id},(1-\alpha) \operatorname{Id}+\alpha \mathrm{P}_{U_{1}}, \ldots,(1-\alpha) \operatorname{Id}+\alpha \mathrm{P}_{U_{m}}\right\}
$$

and that

$$
\mathcal{S}=\left\{\operatorname{Id}, \mathrm{R}_{U_{1}}, \ldots, \mathrm{R}_{U_{m}}\right\}
$$

Then $C C_{\widehat{\mathcal{S}}}$ is proper. Moreover,

$$
\begin{equation*}
(\forall x \in \mathcal{H}) \quad C C_{\widehat{\mathcal{S}}} x=\frac{\alpha}{2} C C_{\mathcal{S}} x+\left(1-\frac{\alpha}{2}\right) x \in \mathcal{H} \tag{5.1}
\end{equation*}
$$

Proof. Apply Proposition 4.12 with $\alpha$ replaced by $\frac{\alpha}{2}$.
Taking $\alpha=1$ in Proposition 5.1, we deduce the next result.
Corollary 5.2. Assume that $\widehat{\mathcal{S}}=\left\{\mathrm{Id}, \mathrm{P}_{U_{1}}, \ldots, \mathrm{P}_{U_{m-1}}, \mathrm{P}_{U_{m}}\right\}$. Then $C C_{\widehat{\mathcal{S}}}$ is proper, that is for every $x \in \mathcal{H}$, there exists unique $C C_{\widehat{\mathcal{S}}} x \in \mathcal{H}$ satisfying
(i) $C C_{\widehat{\mathcal{S}}}(x) \in \operatorname{aff}\left\{x, \mathrm{P}_{U_{1}}(x), \ldots, \mathrm{P}_{U_{m-1}}(x), \mathrm{P}_{U_{m}}(x)\right\}$
(ii) $\left\|C C_{\widehat{\mathcal{S}}}(x)-x\right\|=\left\|C C_{\widehat{\mathcal{S}}}(x)-\mathrm{P}_{U_{1}}(x)\right\|=\cdots=\left\|C C_{\widehat{\mathcal{S}}}(x)-\mathrm{P}_{U_{m}}(x)\right\|$.

Proposition 5.3. Assume that $U_{2}$ is linear and that $\widehat{\mathcal{S}}=\left\{\mathrm{Id}, \mathrm{P}_{U_{1}}, \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}}\right\}$. Then $C C_{\widehat{\mathcal{S}}}$ is proper.
Proof. Let $x \in \mathcal{H}$. If $\operatorname{card}(\widehat{\mathcal{S}}(x)) \leq 2$, by Proposition $3.3, C C_{\widehat{\mathcal{S}}} x \in \mathcal{H}$. Now assume card $(\widehat{\mathcal{S}}(x))=3$. If $x, \mathrm{P}_{U_{1}} x, \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x$ are affinely independent, by Fact 2.14, $C C_{\widehat{\mathcal{S}}} x \in \mathcal{H}$.

Assume that

$$
\begin{equation*}
x, \mathrm{P}_{U_{1}} x, \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x \text { are affinely dependent. } \tag{5.2}
\end{equation*}
$$

Note that card $(\widehat{\mathcal{S}}(x))=3$ implies that $\mathrm{P}_{U_{1}} x-x \neq 0$; moreover, (5.2) yields that there exists $\alpha \neq 1$ such that

$$
\begin{equation*}
\mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x-x=\alpha\left(\mathrm{P}_{U_{1}} x-x\right) \tag{5.3}
\end{equation*}
$$

Because $U_{2}$ is linear subspace, $\mathrm{P}_{U_{2}}$ is linear. Applying to both sides of (5.3) the projector $\mathrm{P}_{U_{2}}$, we obtain

$$
\begin{aligned}
& \mathrm{P}_{U_{2}} \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x-\mathrm{P}_{U_{2}} x=\alpha\left(\mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x-\mathrm{P}_{U_{2}} x\right) \\
\Longrightarrow & \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x-\mathrm{P}_{U_{2}} x=\alpha\left(\mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x-\mathrm{P}_{U_{2}} x\right) \quad(\text { by Lemma } 2.8(\mathrm{i})) \\
\Longrightarrow & (1-\alpha) \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x=(1-\alpha) \mathrm{P}_{U_{2}} x \\
\Longrightarrow & \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x=\mathrm{P}_{U_{2}} x . \quad(\alpha \neq 1)
\end{aligned}
$$

Combining the implications above with card $(\widehat{\mathcal{S}}(x))=3$ and (5.2), we deduce that $x, \mathrm{P}_{U_{1}} x, \mathrm{P}_{U_{2}} x$ are pairwise distinct and affinely dependent. Applying Corollary 5.2 to $m=2$, we obtain $C C\left(\left\{x, \mathrm{P}_{U_{1}} x, \mathrm{P}_{U_{2}} x\right\}\right) \in \mathcal{H}$. But this contradicts Fact 2.14. Therefore, $\operatorname{dom} C C_{\widehat{\mathcal{S}}}=\mathcal{H}$.
Proposition 5.4. Assume that $U_{2}$ is linear and that $\widehat{\mathcal{S}}=\left\{\mathrm{Id}, \mathrm{P}_{U_{2}}, \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}}\right\}$. Then $C C_{\widehat{\mathcal{S}}}$ is proper.
Proof. Let $x \in \mathcal{H}$. Similarly to the proof in Proposition 5.3, we arrive at a contradiction for the case where $\operatorname{card}(\widehat{\mathcal{S}}(x))=3$ and there exists $\alpha \neq 1$ such that

$$
\begin{equation*}
\mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x-x=\alpha\left(\mathrm{P}_{U_{2}} x-x\right) \tag{5.4}
\end{equation*}
$$

As in the proof of Proposition 5.3, we apply to both sides of (5.4) the projector $\mathrm{P}_{U_{2}}$. Then

$$
\begin{aligned}
& \mathrm{P}_{U_{2}} \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x-\mathrm{P}_{U_{2}} x=\alpha\left(\mathrm{P}_{U_{2}} \mathrm{P}_{U_{2}} x-\mathrm{P}_{U_{2}} x\right) \\
\Longrightarrow & \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x-\mathrm{P}_{U_{2}} x=\alpha\left(\mathrm{P}_{U_{2}} x-\mathrm{P}_{U_{2}} x\right)=0 \quad \text { (by Lemma 2.8(i)) } \\
\Longrightarrow & \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x=\mathrm{P}_{U_{2}} x
\end{aligned}
$$

which contradicts card $(\widehat{\mathcal{S}}(x))=3$.
5.2. Improper circumcenter mappings induced by projectors. In view of Propositions 5.1, 5.3 and 5.4, we consider the following question:
Question 5.5. Suppose that $\left\{\alpha_{1}, \alpha_{2}\right\} \subseteq \mathbb{R} \backslash\{0,1\}$ and that at least one of $\alpha_{1}, \alpha_{2}$ is not $2 .{ }^{2}$ Assume that

$$
\widehat{\mathcal{S}}=\left\{\operatorname{Id},\left(1-\alpha_{1}\right) \operatorname{Id}+\alpha_{1} \mathrm{P}_{U_{1}},\left(\left(1-\alpha_{2}\right) \operatorname{Id}+\alpha_{2} \mathrm{P}_{U_{2}}\right)\left(\left(1-\alpha_{1}\right) \operatorname{Id}+\alpha_{1} \mathrm{P}_{U_{1}}\right)\right\}
$$

or that

$$
\widehat{\mathcal{S}}=\left\{\operatorname{Id},\left(1-\alpha_{1}\right) \operatorname{Id}+\alpha_{1} \mathrm{P}_{U_{2}},\left(\left(1-\alpha_{2}\right) \operatorname{Id}+\alpha_{2} \mathrm{P}_{U_{2}}\right)\left(\left(1-\alpha_{1}\right) \operatorname{Id}+\alpha_{1} \mathrm{P}_{U_{1}}\right)\right\}
$$

Is $C C_{\widehat{\mathcal{S}}}$ proper?
The following example demonstrates that the answer to Question 5.5 is negative.
Example 5.6. Assume that $m=2$ and that $U:=U_{1}=U_{2} \varsubsetneqq \mathcal{H}$ and $\left\{\alpha_{1}, \alpha_{2}\right\} \subseteq \mathbb{R}$. Assume further that $\widehat{\mathcal{S}}=\left\{\operatorname{Id},\left(1-\alpha_{1}\right) \operatorname{Id}+\alpha_{1} \mathrm{P}_{U},\left(\left(1-\alpha_{2}\right) \operatorname{Id}+\alpha_{2} \mathrm{P}_{U}\right) \circ((1-\right.$ $\left.\left.\left.\alpha_{1}\right) \operatorname{Id}+\alpha_{1} \mathrm{P}_{U}\right)\right\}$. Then $C C_{\widehat{\mathcal{S}}}$ is improper if and only if $\alpha_{1} \neq 0, \alpha_{1} \neq 1, \alpha_{2} \neq 0$ and $\alpha_{2} \neq \frac{\alpha_{1}}{\alpha_{1}-1}$.
Proof. Since $\mathrm{R}_{U}=2 \mathrm{P}_{U}-\mathrm{Id}$, we deduce that

$$
\begin{aligned}
\widehat{\mathcal{S}} & =\left\{\operatorname{Id},\left(1-\alpha_{1}\right) \mathrm{Id}+\alpha_{1} \mathrm{P}_{U},\left(\left(1-\alpha_{2}\right) \mathrm{Id}+\alpha_{2} \mathrm{P}_{U}\right) \circ\left(\left(1-\alpha_{1}\right) \mathrm{Id}+\alpha_{1} \mathrm{P}_{U}\right)\right\} \\
& =\left\{\mathrm{Id}, \mathrm{Id}+\alpha_{1}\left(\mathrm{P}_{U}-\mathrm{Id}\right),\left(\mathrm{Id}+\alpha_{2}\left(\mathrm{P}_{U}-\mathrm{Id}\right)\right) \circ\left(\mathrm{Id}+\alpha_{1}\left(\mathrm{P}_{U}-\mathrm{Id}\right)\right)\right\} \\
& =\left\{\mathrm{Id}, \mathrm{Id}+\frac{\alpha_{1}}{2}\left(\mathrm{R}_{U}-\mathrm{Id}\right),\left(\mathrm{Id}+\frac{\alpha_{2}}{2}\left(\mathrm{R}_{U}-\mathrm{Id}\right)\right) \circ\left(\mathrm{Id}+\frac{\alpha_{1}}{2}\left(\mathrm{R}_{U}-\mathrm{Id}\right)\right)\right\} .
\end{aligned}
$$

The result now follows from the assumptions above and Example 4.15.

[^2]Next, we present further improper instances of $C C_{\widehat{\mathcal{S}}}$, where Id $\in \widehat{\mathcal{S}} \subseteq$ aff $\Theta$.
Example 5.7. Assume that $\mathcal{H}=\mathbb{R}^{2}, m=2, U_{1}=\mathbb{R} \cdot(1,0)$, and that $U_{2}=\mathbb{R} \cdot(1,2)$. Assume further that $\widehat{\mathcal{S}}=\left\{\mathrm{Id}, \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}}, \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}}\right\}$. Take $x=(2,4) \in U_{2}$. As Figure 6 illustrates, $x, \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x$, and $\mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x$ are pairwise distinct and colinear. By Theorem 3.9, $C C_{\widehat{\mathcal{S}}}$ is improper.


Figure 6. Example 5.7 illustrates $C C_{\widehat{\mathcal{S}}} x=\varnothing$ for the colinear case.
Example 5.8. Assume that $\mathcal{H}=\mathbb{R}^{2}$, that $m=2$, that $U_{1}=\mathbb{R} \cdot(1,0)$, and that $U_{2}=\mathbb{R} \cdot(1,1)$. Assume further that $\widehat{\mathcal{S}}=\left\{\mathrm{Id}, \mathrm{P}_{U_{1}}, \mathrm{P}_{U_{2}}, \mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}}\right\}$. Take $x=(4,2)$ and set $\mathcal{K}=\left\{\mathrm{Id}, \mathrm{P}_{U_{1}}, \mathrm{P}_{U_{2}}\right\}$. Clearly, $\mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x-x \in \mathbb{R}^{2}=\operatorname{span}\left\{\mathrm{P}_{U_{1}} x-x, \mathrm{P}_{U_{2}} x-x\right\}$, which implies that aff $(\mathcal{K}(x))=\operatorname{aff}(\widehat{\mathcal{S}}(x))$. By Fact 2.19, if $C C_{\widehat{\mathcal{S}}} x \in \mathcal{H}$, then $C C_{\widehat{\mathcal{S}}} x=$ $C C_{\mathcal{K}} x$. As Figure 7 shows, $\left\|C C_{\mathcal{K}} x-x\right\|=\left\|C C_{\mathcal{K}} x-\mathrm{P}_{U_{1}} x\right\|=\left\|C C_{\mathcal{K}} x-\mathrm{P}_{U_{2}} x\right\| \neq$ $\left\|C C_{\mathcal{K}} x-\mathrm{P}_{U_{2}} \mathrm{P}_{U_{1}} x\right\|$. Hence $C C_{\widehat{\mathcal{S}}} x=\varnothing$, which implies that $C C_{\widehat{\mathcal{S}}}$ is improper.


Figure 7. Example 5.8 illustrates $C C_{\widehat{\mathcal{S}}} x=\varnothing$ for the non-colinear case.

## 6. More improper circumcenter mappings induced by reflectors

In Theorem 4.3(i), to prove $C C_{\mathcal{S}}$ is proper, we required that
(6.1) $\quad U_{1}, \ldots, U_{m}$ are closed affine subspaces in $\mathcal{H}$ with $\cap_{i=1}^{m} U_{i} \neq \varnothing$,
and that

$$
\begin{equation*}
\operatorname{Id} \in \mathcal{S} \subseteq \Omega \tag{6.2}
\end{equation*}
$$

In Section 4.2, we have already seen that when the condition $\mathcal{S} \subseteq \Omega$ fails, the circumcenter mapping induced by reflectors $C C_{\mathcal{S}}$ may be improper. In the remaining part of this section, we consider two circumcenter mappings induced by reflectors, where $m=2$ and $\mathcal{S}=\left\{\operatorname{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}}\right\}$ or $\mathcal{S}=\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}\right\}$. We construct additional improper circumcenter mappings with the conditions in (6.1) not being satisfied, which means that the conditions (6.1) and (6.2) are sharp.
6.1. Inconsistent cases. In this subsection, we focus on the case when $\cap_{i=1}^{m} U_{i}=$ $\varnothing$. Let $U$ and $V$ be two nonempty, closed, convex (possibly nonintersecting) subsets of $\mathcal{H}$. A best approximation pair relative to $(U, V)$ is

$$
(a, b) \in U \times V \quad \text { such that } \quad\|a-b\|=\inf \|U-V\| .
$$

In the reference [3], the authors used the Douglas-Rachford splitting operator $T=$ $\frac{\mathrm{R}_{V} \mathrm{R}_{U}+\mathrm{Id}}{2}$ to find a best approximation pair relative to $(U, V)$.

Fact 6.1 ([3, Theorem 3.13 and Remark 3.14(ii)]). Let $U$ be a closed affine subspace and let $V$ be a nonempty, closed, convex set in $\mathcal{H}$ ( $U, V$ are possibly nonintersecting). Suppose that best approximation pairs relative to (U,V) exist. Set $T:=\frac{\mathrm{R}_{V} \mathrm{R}_{U}+\mathrm{Id}}{2}$. Let $x_{0} \in \mathcal{H}$ and set $x_{n}=T^{n} x_{0}$, for all $n \in \mathbb{N}$. Then

$$
\left(\left(\mathrm{P}_{V} \mathrm{R}_{U} x_{n}, \mathrm{P}_{U} x_{n}\right)\right)_{n \in \mathbb{N}} \quad \text { and } \quad\left(\left(\mathrm{P}_{V} \mathrm{P}_{U} x_{n}, \mathrm{P}_{U} x_{n}\right)\right)_{n \in \mathbb{N}}
$$

both converge weakly to best approximation pairs relative to $(U, V)$.
The following examples show that even if both of $U_{1}, U_{2}$ are closed affine subspaces, when $U_{1} \cap U_{2}=\varnothing$, the operator $C C_{\mathcal{S}}$ may not be proper where $\mathcal{S}=$ $\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}}\right\}$ or $\mathcal{S}=\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}\right\}$. (Notice that in Example 6.2, $U_{1}$ is even a compact set.) Hence, we can not directly generalize Fact 6.1 by the circumcenter mapping induced by reflectors.

The results of the following examples in this section are easily from Corollary 3.4 and the proofs are omitted.

Example 6.2. Assume that $\mathcal{H}=\mathbb{R}^{2}$, that $U_{1}=\{(2,0)\}$, and that $U_{2}=\mathbb{R} \cdot(0,1)$. Set $\mathcal{S}_{1}=\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}}\right\}$ and $\mathcal{S}_{2}=\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}\right\}$. Then

$$
\begin{aligned}
\operatorname{dom} C C_{\mathcal{S}_{1}} & =\left(\mathbb{R}^{2} \backslash \mathbb{R} \cdot(1,0)\right) \cup\{(2,0),(0,0)\} \\
\operatorname{dom} C C_{\mathcal{S}_{2}} & =\left(\mathbb{R}^{2} \backslash \mathbb{R} \cdot(1,0)\right) \cup\{(2,0),(4,0)\}
\end{aligned}
$$

6.2. Non-affine cases. One of the nice aspects of the Douglas-Rachford method is that it can be used for general convex sets. In this subsection, we assume that

$$
\mathcal{S}_{1}=\left\{\operatorname{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}}\right\} \quad \text { or } \quad \mathcal{S}_{2}=\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}\right\}
$$

We shall present examples in which the operator $C C_{\mathcal{S}}$ is improper, with at least one of $U_{1}$ and $U_{2}$ not being an affine subspace while $U_{1} \cap U_{2} \neq \varnothing$.

Example 6.3. Assume that $\mathcal{H}=\mathbb{R}^{2}$, that $U_{1}=\mathbb{R}_{+}^{2}$, and that $U_{2}=(2,0)+\mathbb{R} \cdot(0,1)$. Then

$$
\begin{aligned}
\operatorname{dom} C C_{\mathcal{S}_{1}} & =\mathbb{R}^{2} \backslash\{(x, y) \mid x<0 \text { and } y \geq 0\} \\
\operatorname{dom} C C_{\mathcal{S}_{2}} & =\left(\mathbb{R}^{2} \backslash\{(x, y) \mid x<0 \text { and } y \geq 0\}\right) \cup\{(-2, y) \mid y \geq 0\}
\end{aligned}
$$

In the remainder of this subsection, we revisit the examples used in [6] to show the potential of the Circumcentering Douglas-Rachford method, which are the iterations of the operator $C C_{\mathcal{S}_{2}}$.

Example 6.4. Assume that $\mathcal{H}=\mathbb{R}^{2}$, that $U_{1}=\mathbf{B}[(0,0) ; 1]$, and that $U_{2}=(1,0)+$ $\mathbb{R} \cdot(0,1)$. Then

$$
\begin{aligned}
\operatorname{dom} C C_{\mathcal{S}_{1}} & =\mathbb{R}^{2} \backslash\{(x, 0) \mid x<-1\} \\
\operatorname{dom} C C_{\mathcal{S}_{2}} & =\mathbb{R}^{2} \backslash\{(x, 0) \mid x<-3 \text { or }-3<x<-1\}
\end{aligned}
$$

Example 6.5. Assume that $\mathcal{H}=\mathbb{R}^{2}$, that $U_{1}=\mathbf{B}[(0,0) ; 1]$, and that $U_{2}=\mathbb{R} \cdot(0,1)$. Then

$$
\begin{aligned}
\operatorname{dom} C C_{\mathcal{S}_{1}} & =\mathbb{R}^{2} \backslash\{(x, 0)| | x \mid>1\} \\
\operatorname{dom} C C_{\mathcal{S}_{2}} & =\mathbb{R}^{2} \backslash\{(x, 0)| | x \mid>2 \text { or } 1<|x|<2\}
\end{aligned}
$$

Example 6.6. Assume that $\mathcal{H}=\mathbb{R}^{2}$, that $U_{1}=\mathbf{B}[(-1,0) ; 1]$, and that $U_{2}=$ $\mathbf{B}[(1,0) ; 1]$. Then

$$
\begin{gathered}
\operatorname{dom} C C_{\mathcal{S}_{1}}=\mathbb{R}^{2} \backslash\{(x, 0) \mid x<-2 \text { or } x>2\}, \\
\{(x, 0) \mid-6 \leq x \leq-4 \text { or } x \geq-2\} \subseteq \operatorname{dom} C C_{\mathcal{S}_{2}} \\
\{(x, 0) \mid x<-6 \text { or }-4<x<-2\} \subseteq \mathbb{R}^{2} \backslash\left(\operatorname{dom} C C_{\mathcal{S}_{2}}\right)
\end{gathered}
$$

Example 6.7. Assume that $\mathcal{H}=\mathbb{R}^{2}$, that $U_{1}=\mathbf{B}[(-1,0) ; 2]$, and that $U_{2}=$ $\mathbf{B}[(1,0) ; 2]$. Then

$$
\begin{gathered}
\operatorname{dom} C C_{\mathcal{S}_{1}}=\mathbb{R}^{2} \backslash\{(x, 0) \mid x<-3 \text { or } x>3\} \\
\{(x, 0) \mid-9 \leq x \leq-5 \text { or }-3 \leq x \leq 3\} \subseteq \operatorname{dom} C C_{\mathcal{S}_{2}} \\
\{(x, 0) \mid x<-9 \text { or }-5<x<-3 \text { or } x>3\} \subseteq \mathbb{R}^{2} \backslash\left(\operatorname{dom} C C_{\mathcal{S}_{2}}\right) .
\end{gathered}
$$

Finally, consider $U_{1}=\left\{(x, y) \in \mathbb{R}^{2} \mid(x+1)^{2}+y^{2}=4\right\}$ and $U_{2}=\left\{(x, y) \in \mathbb{R}^{2} \mid(x-\right.$ $\left.1)^{2}+y^{2}=4\right\}$. Note that neither $U_{1}$ nor $U_{2}$ is convex. For $\mathcal{S}=\left\{\operatorname{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}\right\}$ or $\mathcal{S}=\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}}\right\}$, one can show that $\operatorname{dom} C C_{\mathcal{S}} \varsubsetneqq \mathbb{R}^{2}$.
6.3. Impossibility to extend to maximally monotone operators. Assume that $\mathcal{S}=\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}}\right\}$ or $\mathcal{S}=\left\{\mathrm{Id}, \mathrm{R}_{U_{1}}, \mathrm{R}_{U_{2}} \mathrm{R}_{U_{1}}\right\}$. In order to show a counterexample where the definition of $C C_{\mathcal{S}}$ fails to be directly generalized to maximally monotone theory, we need the definition and facts below.

Definition 6.8 ([2, Definition 23.1]). Let $A: \mathcal{H} \rightarrow 2^{\mathcal{H}}$. The resolvent of $A$ is

$$
J_{A}=(\operatorname{Id}+A)^{-1}
$$

Fact 6.9 ([2, Corollary 23.11]). Let $A: \mathcal{H} \rightarrow 2^{\mathcal{H}}$ be maximally monotone and let $\gamma \in \mathbb{R}_{++}$. Then the following hold.
(i) $J_{\gamma A}: \mathcal{H} \rightarrow \mathcal{H}$ and $\operatorname{Id}-J_{\gamma A}: \mathcal{H} \rightarrow \mathcal{H}$ are firmly nonexpansive and maximally monotone.
(ii) The reflected resolvent

$$
R_{\gamma A}: \mathcal{H} \rightarrow \mathcal{H}: x \mapsto 2 J_{\gamma A} x-x
$$

is nonexpansive.
Fact 6.10 ([2, Corollary 20.28]). Let $A: \mathcal{H} \rightarrow \mathcal{H}$ be monotone and continuous. Then $A$ is maximally monotone.

Fact 6.11 ([2, Corollary 20.26]). Let $C$ be a nonempty closed convex subset of $\mathcal{H}$. Then $N_{C}$ is maximally monotone.

Fact 6.12 ([2, Corollary 23.4]). Let $C$ be a nonempty closed convex subset of $\mathcal{H}$. Then

$$
J_{N_{C}}=\left(\operatorname{Id}+N_{C}\right)^{-1}=\mathrm{P}_{C}
$$

By Fact 6.12, $\mathrm{R}_{U_{1}}=2 \mathrm{P}_{U_{1}}-\mathrm{Id}=2 J_{N_{U_{1}}}-\mathrm{Id}$ and $\mathrm{R}_{U_{2}}=2 \mathrm{P}_{U_{2}}-\mathrm{Id}=2 J_{N_{U_{2}}}-\mathrm{Id}$. In these special cases, the reflectors are consistent with the corresponding reflected resolvent.

In the following examples, we replace the two maximally monotone operators $N_{U_{1}}, N_{U_{2}}$ in the set $\mathcal{S}=\left\{\operatorname{Id}, 2 J_{N_{U_{1}}}-\mathrm{Id}, 2 J_{N_{U_{2}}}-\mathrm{Id}\right\}$ or $\mathcal{S}=\left\{\operatorname{Id}, 2 J_{N_{U_{1}}}-\mathrm{Id},\left(2 J_{N_{U_{2}}}-\right.\right.$ Id) $\left.\circ\left(2 J_{N_{U_{1}}}-\mathrm{Id}\right)\right\}$ by $\alpha$ Id and $\beta$ Id respectively, with $\alpha \geq 0$ and $\beta \geq 0$. By Fact 6.10, since $\alpha \geq 0$ and $\beta \geq 0$, we obtain that $\alpha$ Id and $\beta$ Id are maximally monotone operators. We shall characterize the improperness of the new operator $C C_{\mathcal{S}}$.

Example 6.13. Assume that $\{0\} \varsubsetneqq \mathcal{H}$. Set $A=\alpha \mathrm{Id}$ and $B=\beta \mathrm{Id}$, where $\alpha \geq 0$ and $\beta \geq 0$. Further set

$$
\mathcal{S}_{1}=\left\{\operatorname{Id}, R_{A}, R_{B}\right\} \quad \text { and } \quad \mathcal{S}_{2}=\left\{\operatorname{Id}, R_{A}, R_{B} R_{A}\right\}
$$

Then $C C_{\mathcal{S}_{1}}$ is improper if and only if $\alpha \neq 0, \beta \neq 0$ and $\alpha \neq \beta$. Moreover, $C C_{\mathcal{S}_{2}}$ is improper if and only if $\alpha \neq 0, \alpha \neq 1, \beta \neq 0$ and $\alpha \neq-\beta$.

Proof. The definitions yield

$$
\begin{aligned}
& J_{A}=(A+\mathrm{Id})^{-1}=((\alpha+1) \mathrm{Id})^{-1}=\frac{1}{\alpha+1} \mathrm{Id} \\
& R_{A}=2 J_{A}-\mathrm{Id}=\frac{2}{\alpha+1} \mathrm{Id}-\mathrm{Id}=\frac{1-\alpha}{\alpha+1} \mathrm{Id} \\
& J_{B}=(B+\mathrm{Id})^{-1}=\left(((\beta+1) \mathrm{Id})^{-1}=\frac{1}{\beta+1} \mathrm{Id}\right. \\
& R_{B}=2 J_{B}-\mathrm{Id}=\frac{2}{\beta+1} \mathrm{Id}-\mathrm{Id}=\frac{1-\beta}{\beta+1} \mathrm{Id}
\end{aligned}
$$

Let $x \in \mathcal{H} \backslash 0$. Now

$$
\begin{equation*}
x=R_{A} x \Longleftrightarrow x=\frac{1-\alpha}{\alpha+1} x \Longleftrightarrow 1=\frac{1-\alpha}{\alpha+1} \Longleftrightarrow \alpha=0 \tag{6.3a}
\end{equation*}
$$

$$
\begin{equation*}
x=R_{B} x \Longleftrightarrow x=\frac{1-\beta}{\beta+1} x \Longleftrightarrow \beta=0 \tag{6.3b}
\end{equation*}
$$

$$
\begin{equation*}
R_{A} x=R_{B} x \Longleftrightarrow \frac{1-\alpha}{\alpha+1} x=\frac{1-\beta}{\beta+1} x \Longleftrightarrow \alpha=\beta \tag{6.3c}
\end{equation*}
$$

$$
\begin{equation*}
x=R_{B} R_{A} x \Longleftrightarrow x=\frac{1-\alpha}{\alpha+1} \frac{1-\beta}{\beta+1} x \Longleftrightarrow \alpha=-\beta \tag{6.3~d}
\end{equation*}
$$

$$
\begin{equation*}
R_{A} x=R_{B} R_{A} x \Longleftrightarrow \frac{1-\alpha}{\alpha+1} x=\frac{1-\alpha}{\alpha+1} \frac{1-\beta}{\beta+1} x \Longleftrightarrow \alpha=1 \text { or } \beta=0 \tag{6.3e}
\end{equation*}
$$

" $\Longrightarrow$ ": According to the previous analysis, in both of the assertions, the contrapositive of the required results follow from Proposition 3.3.
$" \Longleftarrow ":$ Assume $\alpha \neq 0, \beta \neq 0$ and $\alpha \neq \beta$. Then

$$
\begin{aligned}
\operatorname{aff}\left(\mathcal{S}_{1}(x)\right) & =\operatorname{aff}\left\{x, R_{A} x, R_{B} x\right\}=x+\operatorname{span}\left\{R_{A} x-x, R_{B} x-x\right\} \\
& =x+\operatorname{span}\left\{\frac{-2 \alpha}{\alpha+1} x, \frac{-2 \beta}{\beta+1} x\right\} \\
& =\mathbb{R} \cdot x
\end{aligned}
$$

Let $x \in \mathcal{H} \backslash\{0\}$. We observe that

$$
\begin{aligned}
& \left(\exists y \in \operatorname{aff}\left(\mathcal{S}_{1}(x)\right)\right) \quad\|y-x\|=\left\|y-R_{A} x\right\|=\left\|y-R_{B} x\right\| \\
\Longleftrightarrow & (\exists t \in \mathbb{R}) \quad\|t x-x\|=\left\|t x-\frac{1-\alpha}{\alpha+1} x\right\|=\left\|t x-\frac{1-\beta}{\beta+1} x\right\| \\
\Longleftrightarrow & (\exists t \in \mathbb{R}) \quad|t-1|=\left|t-\frac{1-\alpha}{\alpha+1}\right|=\left|t-\frac{1-\beta}{\beta+1}\right| . \quad(\text { by } x \neq 0)
\end{aligned}
$$

On the other hand, combining the assumptions with Corollary 2.16 and (6.3), we obtain that

$$
(\nexists t \in \mathbb{R}) \quad|t-1|=\left|t-\frac{1-\alpha}{\alpha+1}\right|=\left|t-\frac{1-\beta}{\beta+1}\right|
$$

Hence,

$$
(\forall x \in \mathcal{H} \backslash\{0\}) \quad C C_{\mathcal{S}_{1}} x=\varnothing
$$

Assume $\alpha \neq 0, \alpha \neq 1, \beta \neq 0$ and $\alpha \neq-\beta$. A similar proof shows that for every $x \in \mathcal{H} \backslash\{0\}$, there is no point $y \in \operatorname{aff}\left(\mathcal{S}_{2}(x)\right)$, such that $\|y-x\|=\left\|y-R_{A} x\right\|=$ $\left\|y-R_{B} R_{A} x\right\|$, which implies that $(\forall x \in \mathcal{H} \backslash\{0\}) C C_{\mathcal{S}_{2}} x=\varnothing$.

Arguing similarly to the proof of the previous result, we also obtain the following result:

Example 6.14. Assume that $\{0\} \varsubsetneqq \mathcal{H}$. Let $\{a, b\} \subseteq \mathbb{R}$. Set $A \equiv a$, i.e., $(\forall x \in \mathcal{H})$ $A x=a$, and $B \equiv b$. Furthermore, set

$$
\mathcal{S}_{1}=\left\{\operatorname{Id}, R_{A}, R_{B}\right\} \quad \text { and } \quad \mathcal{S}_{2}=\left\{\operatorname{Id}, R_{A}, R_{B} R_{A}\right\}
$$

Then $C C_{\mathcal{S}_{1}}$ is improper if and only if $a \neq 0, b \neq 0$ and $a \neq b$. Moreover, $C C_{\mathcal{S}_{2}}$ is improper if and only if $a \neq 0, b \neq 0$ and $a \neq-b$.

The example above shows that there is no direct way to generalize the definition of $C C_{\mathcal{S}}$ to maximally monotone theory.

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[^1]:    ${ }^{1}$ When $\mathcal{S}(x)$ is a singleton, then $\widetilde{\mathcal{S}}_{x}=\left\{T_{1}\right\}$ by the standard convention that $\varnothing$ is the basis of $\{0\}$.

[^2]:    ${ }^{2}$ For $i \in\{1,2\}$, when $\alpha_{i}$ is 0,1 , or 2 , then $\left(1-\alpha_{i}\right) \operatorname{Id}+\alpha_{i} \mathrm{P}_{U_{1}}$ is $\mathrm{Id}, \mathrm{P}_{U_{1}}$, or $\mathrm{R}_{U_{1}}$ respectively. In these special cases, the answer for Question 5.5 is positive (see Proposition 5.1 and Theorem 4.3(i)).

