# Self-adaptive CQ-type algorithms for the split feasibility problem involving two bounded linear operators in Hilbert spaces 

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#### Abstract

In this article, we consider and investigate a split convex feasibility problem involving two bounded linear operators in Hilbert spaces. We introduce a self-adaptive CQ-type algorithm by selecting the stepsize which is independent of the operator norms and establish a strong convergence result of the proposed algorithm under some mild control conditions. Moreover, we propose a self-adaptive relaxed CQ-type algorithm for solving the problem constrained by sub-level sets of convex functions. A numerical example and an application in compressed sensing are also given to illustrate the convergence behaviour of our proposed algorithms. Our results in this paper improve and generalize some existing results in the literature.


## 1. Introduction

Let $C$ and $Q$ be two nonempty closed convex subsets of Hilbert spaces $H_{1}$ and $H_{2}$, respectively. The split feasibility problem (shortly, SFP) is to find a point

$$
\begin{equation*}
x \in C \text { such that } A x \in Q \text {, } \tag{1.1}
\end{equation*}
$$

where $A: H_{1} \rightarrow H_{2}$ is a bounded linear operator. The SFP is the first instance of the split inverse problem (referred to [13, Sect. 2]), which was first introduced by Censor and Elfving [11] in Euclidean spaces. The SFP model can be applied to solving many mathematical problems such as the constrained least-squares problem, the linear split feasibility problem, and the linear programming problem and it can be used in real-world applications, for example, in signal processing, in image recovery, in intensity-modulated therapy, in pattern recognition and in data prediction (see $[3,5,10,12,20,22]$ ). Consequently, the SFP has been widely studied and various methods for solving such a problem have been invented and developed by many authors, see [ $2,9,17,24,25,35,36,37,38,41,43,44]$ and the references therein. One of the powerful methods for approximating solutions of (1.1) is known as the CQ algorithm introduced by Byrne [2] as follows:

$$
\left\{\begin{array}{l}
x_{1} \in H_{1}  \tag{1.2}\\
x_{k+1}=P_{C}\left(x_{k}-\lambda A^{*}\left(I-P_{Q}\right) A x_{k}\right), \quad k \geq 1
\end{array}\right.
$$

where $\lambda \in\left(0,2 /\|A\|^{2}\right), P_{C}$ and $P_{Q}$ are the metric projections onto $C$ and $Q$, respectively, and $A^{*}$ stands for the adjoint operator of $A$. After that, various kinds of the split inverse problem, which are generalizations of the SFP were introduced and studied, see [4, 12, 13, $14,28,32$ ] for instance.

[^0]In this paper, we focus on a generalization of the SFP (1.1) in which two bounded linear operators $A, B: H_{1} \rightarrow H_{2}$ are involved - Finding a point

$$
\begin{equation*}
x \in C \text { such that } A x \in Q \text { and } B x \in \tilde{Q} \tag{1.3}
\end{equation*}
$$

where $C \subseteq H_{1}$ and $Q, \tilde{Q} \subseteq H_{2}$ are nonempty closed convex subsets. We call (1.3) the two-operator split feasibility problem (two-operator SFP), see [28, 32] for the general versions of this problem. The two-operator SFP (1.3) can be reduced to the convexly constrained linear problem $([15,29])$ involving two linear operators, that is, finding a point $x \in C$ such that $A x=y, B x=\tilde{y}$ in $H_{2}$.

In 2019, Kangtunyakarn [21] studied the two-operator SFP (1.3) in case that $Q=\tilde{Q}$ and introduced a viscosity-based algorithm with a given contraction $f: C \rightarrow C$ as follows:
$\left\{\begin{array}{l}x_{1} \in C, \\ x_{k+1}=\beta_{k} f\left(x_{k}\right)+\delta_{k} x_{k}+\gamma_{k} P_{C}\left[x_{k}-\frac{\lambda}{2}\left(A^{*}\left(I-P_{Q}\right) A x_{k}+B^{*}\left(I-P_{Q}\right) B x_{k}\right)\right], \quad k \geq 1,\end{array}\right.$
where $\lambda \in\left(0,2 / \max \left\{\|A\|^{2},\|B\|^{2}\right\}\right)$ and $\left\{\beta_{k}\right\},\left\{\delta_{k}\right\}$, and $\left\{\gamma_{k}\right\}$ are real sequences in $(0,1)$. A strong convergence theorem of (1.4) was proved under some suitable conditions on the control sequences, see [21, Theorem 3.1].

It is noted that the parameters $\lambda$ in (1.2) and in (1.4) depend on the norms of bounded linear operators, so these algorithms have a drawback in the sense that the implementation of them requires to calculate or estimate the operator norms, which is not an easy task in general practice (see [25, Subsection 6.1.2] for instance). To overcome this, in [2, Proposition 4.1], it was presented a helpful method for estimating operator (matrix) norms but its conditions seem restrictive. López et al. [25] proposed an alternative way that is to select the stepsize $\lambda_{k}$ which does not need any prior knowledge of the operator norm for replacing the parameter $\lambda$ in (1.2) as follows:

$$
\begin{equation*}
\lambda_{k}:=\frac{\mu_{k}\left\|\left(I-P_{Q}\right) A x_{k}\right\|^{2}}{2\left\|A^{*}\left(I-P_{Q}\right) A x_{k}\right\|^{2}}, \tag{1.5}
\end{equation*}
$$

where $\mu_{k} \in(0,4)$. We can see that the choice of the stepsize $\lambda_{k}$ in (1.5) is independent of the operator norm $\|A\|$. This stepsize was widely employed in optimization methods and was also modified for use in fixed point methods, see [ $8,18,19,27,33,35]$. The CQ algorithm with the self-adaptive stepsize defined by (1.5) [25, Algorithm 3.1] guarantees only weak convergence for the SFP (1.1), see [25, Theorem 3.5]. However, strong convergence gives more desirable theoretical result in the setting of Hilbert spaces. To get strong convergence, Vinh et al. [35] employed a modification of the CQ algorithm ([37, Algorithm 4.1]) with the stepsize (1.5) for solving the SFP (1.1) as follows:

$$
\left\{\begin{array}{l}
x_{1} \in H_{1}  \tag{1.6}\\
x_{k+1}=P_{C}\left[\left(1-\beta_{k}\right)\left(x_{k}-\lambda_{k} A^{*}\left(I-P_{Q}\right) A x_{k}\right)\right], \quad k \geq 1
\end{array}\right.
$$

where the stepsize $\lambda_{k}$ is defined by (1.5) and $\left\{\beta_{k}\right\} \subset(0,1)$. They proved that the sequence generated by (1.6) converges strongly to the minimum-norm solution to (1.1) under some suitable control conditions, see [35, Theorem 3.1].

Here, the above review leads us to the following natural questions.

1. Can we design a CQ-type algorithm whose stepsize does not depend on the operator norm $\|A\|$ or $\|B\|$ to solve the two-operator SFP (1.3)?
2. How do we adapt the algorithm designed from Question 1 to be a strongly convergent method?

Motivated and inspired by the above questions and the results of Kangtunyakarn [21], López et al. [25], and Vinh et al. [35], we aim to invent a self-adaptive CQ-type algorithm whose stepsize does not depend on any operator norms for solving the two-operator SFP in the setting of Hilbert spaces. Moreover, we will prove that the sequence generated by the proposed algorithm converges strongly to the minimum-norm solution. The rest of the paper is organized as follows. In Sect. 2, some basic facts and useful lemmas for proving our main results are given. Our main result is in Sect. 3. In this section, we introduce a self-adaptive CQ-type algorithm using the stepsize which is independent of the bounded linear operator norms for finding a solution of (1.3). A strong convergence theorem of the proposed algorithm is analyzed and established. In Sect. 4, we propose a self-adaptive relaxed CQ-type algorithm for solving the two-operator SFP in case of sub-level sets of convex functions and also prove its strong convergence result. Finally, in Sect. 5, we provide numerical experiments of our proposed algorithms in the setting of a Euclidean space and in the signal recovery problem with two different blurring operations, and also compare the efficiency of our algorithms with that of some methods depending on the operator norms.

## 2. Preliminaries

Throughout this paper, we suppose that $H, H_{1}$ and $H_{2}$ are real Hilbert spaces with inner products $\langle\cdot, \cdot\rangle$ and the induced norms $\|\cdot\|$ (in particular, in Euclidean spaces, $\|\cdot\|_{1}$ denotes the $l_{1}$-norm and $\|\cdot\|_{2}$ denotes the Euclidean norm). The notation $I$ stands for the identity operator on a Hilbert space. Let $\left\{x_{k}\right\}$ be a sequence in $H$. Weak and strong convergence of $\left\{x_{k}\right\}$ to $x \in H$ are denoted by $x_{k} \rightharpoonup x$ and $x_{k} \rightarrow x$, respectively. The set of all weak-cluster points of $\left\{x_{k}\right\}$ is denoted by $\omega_{w}\left(x_{k}\right)$.

Let $f: H \rightarrow \mathbb{R}$ be a function and $x \in H$. We say that $f$ is weakly lower semi-continuous at $x$ if for every sequence $\left\{x_{k}\right\} \subset H, x_{k} \rightharpoonup x$ implies $f(x) \leq \liminf _{k \rightarrow \infty} f\left(x_{k}\right)$. A subdifferential $\partial f$ of $f$ at $x$ is defined by

$$
\partial f(x)=\{u \in H: f(x)+\langle u, z-x\rangle \leq f(z), \forall z \in H\} .
$$

The function $f$ is said to be subdifferentiable at $x$ if $\partial f(x) \neq \emptyset$. One can see that if $f$ is subdifferentiable at $x$, then $f$ is weakly lower semi-continuous at $x$. We denote the gradient of $f$ by $\nabla f$ if $f$ is differentiable.

Let $K$ be a nonempty closed convex subset of $H$. Recall that the metric projection $P_{K}$ from $H$ onto $K$ assigns to each $x \in H$ the unique point $P_{K} x$ in $K$ satisfying $\left\|x-P_{K} x\right\|=$ $\inf _{z \in K}\|x-z\|$. Some properties of the metric projection are listed below.

Lemma 2.1. The metric projection $P_{K}$ has the following properties:
(1) $\left\langle x-P_{K} x, z-P_{K} x\right\rangle \leq 0, \quad \forall x \in H, \forall z \in K$;
(2) $\left\langle x-P_{K} x, x-z\right\rangle \geq\left\|x-P_{K} x\right\|^{2}, \quad \forall x \in H, \forall z \in K$;
(3) $P_{K}$ is firmly nonexpansive, i.e.,

$$
\left\|P_{K} x-P_{K} y\right\|^{2} \leq\|x-y\|^{2}-\left\|\left(x-P_{K} x\right)-\left(y-P_{K} y\right)\right\|^{2}, \quad \forall x, y \in H
$$

in particular,

$$
\left\|P_{K} x-z\right\|^{2} \leq\|x-z\|^{2}-\left\|x-P_{K} x\right\|^{2}, \quad \forall x \in H, \forall z \in K
$$

Let $Q$ be a nonempty closed convex subset of $H_{2}$ and let $A: H_{1} \rightarrow H_{2}$ be a bounded linear operator with the adjoint operator $A^{*}$. Define a function $f: H_{1} \rightarrow \mathbb{R}$ by

$$
f(x)=\frac{1}{2}\left\|\left(I-P_{Q}\right) A x\right\|^{2} .
$$

We know that $f$ is weakly lower semi-continuous on $H_{1}$ and differentiable with the gradient $\nabla f: H_{1} \rightarrow H_{1}$ given by

$$
\nabla f(x)=A^{*}\left(I-P_{Q}\right) A x
$$

Moreover, $\nabla f$ is Lipschitz continuous with the Lipschitz constant $\|A\|^{2}$, i.e.,

$$
\|\nabla f(x)-\nabla f(y)\| \leq\|A\|^{2}\|x-y\|, \quad \forall x, y \in H_{1}
$$

For more details, the reader is referred to optimization books, see $[1,31]$ for instance.
We end this section with the following useful lemmas for proving our strong convergence results.

Lemma 2.2 ([39]). Let $\left\{t_{k}\right\}$ be a sequence of nonnegative real numbers satisfying

$$
t_{k+1} \leq\left(1-\beta_{k}\right) t_{k}+\beta_{k} \delta_{k}, \quad \forall k \in \mathbb{N}
$$

where $\left\{\beta_{k}\right\}$ is a sequence in $(0,1)$ and $\left\{\delta_{k}\right\}$ is a sequence of real numbers such that $\sum_{k=1}^{\infty} \beta_{k}=\infty$ and $\limsup _{k \rightarrow \infty} \delta_{k} \leq 0$. Then, $\lim _{k \rightarrow \infty} t_{k}=0$.

Lemma 2.3 ([26]). Let $\left\{s_{k}\right\}$ be a sequence of real numbers such that there exists a subsequence $\left\{k_{j}\right\}$ of $\{k\}$ which satisfies $s_{k_{j}}<s_{k_{j}+1}$ for all $j \in \mathbb{N}$. Let $\{\tau(k)\}$ be a sequence of positive integers defined by

$$
\tau(k):=\max \left\{n \leq k: s_{n}<s_{n+1}\right\}
$$

for all $k \geq k_{0}$ (for some $k_{0}$ large enough). Then $\{\tau(k)\}$ is a nondecreasing sequence such that $\tau(k) \rightarrow \infty$ as $k \rightarrow \infty$, and it holds that

$$
s_{\tau(k)} \leq s_{\tau(k)+1} \quad \text { and } s_{k} \leq s_{\tau(k)+1}, \quad \forall k \geq k_{0}
$$

## 3. Self-ADAptive CQ-Type algorithm and its convergence result

This main section provides positive answers to the questions raised in the introduction section, namely that we introduce a self-adaptive CQ-type algorithm (Algorithm 1) whose stepsize does not depend on the bounded linear operator norms to solve the two-operator SFP (1.3). Subsequently, we analyze and establish strong convergence of the proposed algorithm. The following assumptions are set throughout the section:

- $H_{1}$ and $H_{2}$ are real Hilbert spaces,
- $C \subseteq H_{1}$ and $Q, \tilde{Q} \subseteq H_{2}$ are nonempty closed convex subsets,
- $A, B: H_{1} \rightarrow H_{2}$ are two bounded linear operators,
- $\Gamma=\{x \in C: A x \in Q, B x \in \tilde{Q}\} \neq \emptyset$.

We begin with the following result that will be helpful in designing our algorithm.
Lemma 3.4. Let $x^{*} \in C$. Then, $x^{*} \in \Gamma$ if and only if $\left\|A^{*}\left(I-P_{Q}\right) A x^{*}+B^{*}\left(I-P_{\tilde{Q}}\right) B x^{*}\right\|=$ 0 .

Proof. Let $x^{*} \in C$. If $x^{*} \in \Gamma$, then $A x^{*} \in Q, B x^{*} \in \tilde{Q}$ and so $\left(I-P_{Q}\right) A x^{*}=\left(I-P_{\tilde{Q}}\right) B x^{*}=$ 0 . It is obvious that $\left\|A^{*}\left(I-P_{Q}\right) A x^{*}+B^{*}\left(I-P_{\tilde{Q}}\right) B x^{*}\right\|=0$. Conversely, we assume that

$$
\left\|A^{*}\left(I-P_{Q}\right) A x^{*}+B^{*}\left(I-P_{\tilde{Q}}\right) B x^{*}\right\|=0 . \text { Pick } p \in \Gamma . \text { By Lemma 2.1(2), we have }
$$

$$
\begin{aligned}
0 & =\left\|A^{*}\left(I-P_{Q}\right) A x^{*}+B^{*}\left(I-P_{\tilde{Q}}\right) B x^{*}\right\|\left\|x^{*}-p\right\| \\
& \geq\left\langle A^{*}\left(I-P_{Q}\right) A x^{*}+B^{*}\left(I-P_{\tilde{Q}}\right) B x^{*}, x^{*}-p\right\rangle \\
& =\left\langle A^{*}\left(I-P_{Q}\right) A x^{*}, x^{*}-p\right\rangle+\left\langle B^{*}\left(I-P_{\tilde{Q}}\right) B x^{*}, x^{*}-p\right\rangle \\
& =\left\langle\left(I-P_{Q}\right) A x^{*}, A x^{*}-A p\right\rangle+\left\langle\left(I-P_{\tilde{Q}}\right) B x^{*}, B x^{*}-B p\right\rangle \\
& \geq\left\|\left(I-P_{Q}\right) A x^{*}\right\|^{2}+\left\|\left(I-P_{\tilde{Q}}\right) B x^{*}\right\|^{2},
\end{aligned}
$$

which implies that $\left(I-P_{Q}\right) A x^{*}=\left(I-P_{\tilde{Q}}\right) B x^{*}=0$. Hence, $A x^{*} \in Q$ and $B x^{*} \in \tilde{Q}$, that is, $x^{*} \in \Gamma$.

Here, our iterative algorithm for solving the two-operator SFP (1.3) is designed as follows.

```
Algorithm 1: Self-adaptive CQ-type algorithm for the two-operator SFP
    Initialization: Take two real sequences \(\left\{\beta_{k}\right\} \subset(0,1)\) and \(\left\{\mu_{k}\right\} \subset(0,4)\).
    Choose \(x_{0} \in H_{1}\) arbitrarily. Set \(x_{1}=P_{C} x_{0}\) and \(k=1\).
    Iterative Step: Given \(x_{k}\), if \(\left\|A^{*}\left(I-P_{Q}\right) A x_{k}+B^{*}\left(I-P_{\tilde{Q}}\right) B x_{k}\right\|=0\), then
        \(x_{k+1}=x_{k}\) (in this case, \(x_{k}\) solves (1.3) by Lemma 3.4) and the iterative process
        stops. Otherwise, calculate
\[
\begin{gather*}
\lambda_{k}=\mu_{k} \frac{\left\|\left(I-P_{Q}\right) A x_{k}\right\|^{2}+\left\|\left(I-P_{\tilde{Q}}\right) B x_{k}\right\|^{2}}{\left\|A^{*}\left(I-P_{Q}\right) A x_{k}+B^{*}\left(I-P_{\tilde{Q}}\right) B x_{k}\right\|^{2}}  \tag{3.7}\\
x_{k+1}=P_{C}\left[\left(1-\beta_{k}\right)\left(x_{k}-\frac{\lambda_{k}}{2}\left(A^{*}\left(I-P_{Q}\right) A x_{k}+B^{*}\left(I-P_{\tilde{Q}}\right) B x_{k}\right)\right)\right] . \tag{3.8}
\end{gather*}
\]
```

Update $k:=k+1$ and return to Iterative Step.

For the sake of simplicity, we let $g: H_{1} \rightarrow \mathbb{R}$ be defined by

$$
\begin{equation*}
g(x):=\frac{1}{4}\left(\left\|\left(I-P_{Q}\right) A x\right\|^{2}+\left\|\left(I-P_{\tilde{Q}}\right) B x\right\|^{2}\right) \tag{3.9}
\end{equation*}
$$

with the gradient given by

$$
\nabla g(x)=\frac{1}{2}\left(A^{*}\left(I-P_{Q}\right) A x+B^{*}\left(I-P_{\tilde{Q}}\right) B x\right), \quad x \in H_{1}
$$

Note that (1.3) is equivalent to the problem of finding $x \in C$ such that $g(x)=0$. In other words, (3.7) and (3.8) can be rewritten in the form of the following modified gradientprojection method:

$$
x_{k+1}=P_{C}\left[\left(1-\beta_{k}\right)\left(x_{k}-\lambda_{k} \nabla g\left(x_{k}\right)\right)\right], \text { where } \lambda_{k}=\frac{\mu_{k} g\left(x_{k}\right)}{\left\|\nabla g\left(x_{k}\right)\right\|^{2}} .
$$

To verify the convergence of Algorithm 1, the following two lemmas are required.

Lemma 3.5. Let $\left\{x_{k}\right\}$ be a sequence generated by Algorithm 1. If $\nabla g\left(x_{k}\right) \neq 0$, then the following two inequalities hold for all $x^{*} \in \Gamma$,

$$
\begin{equation*}
\left\|x_{k+1}-x^{*}\right\|^{2} \leq \beta_{k}\left\|x^{*}\right\|^{2}+\left(1-\beta_{k}\right)\left\|x_{k}-x^{*}\right\|^{2}-\mu_{k}\left(4-\mu_{k}\right)\left(1-\beta_{k}\right) \frac{g^{2}\left(x_{k}\right)}{\left\|\nabla g\left(x_{k}\right)\right\|^{2}} \tag{3.10}
\end{equation*}
$$

and

$$
\left\|x_{k+1}-x^{*}\right\|^{2} \leq\left(1-\beta_{k}\right)\left\|x_{k}-x^{*}\right\|^{2}+\beta_{k}\left[\beta_{k}\left\|x^{*}\right\|^{2}+2\left(1-\beta_{k}\right)\left\langle x_{k}-x^{*},-x^{*}\right\rangle\right.
$$

$$
\begin{equation*}
\left.+2\left(1-\beta_{k}\right) \lambda_{k}\left\langle\nabla g\left(x_{k}\right), x^{*}\right\rangle\right] \tag{3.11}
\end{equation*}
$$

Proof. Let $x^{*} \in \Gamma$. Using (3.8) and Lemma 2.1(3), we have

$$
\begin{align*}
\left\|x_{k+1}-x^{*}\right\|^{2} & =\left\|P_{C}\left[\left(1-\beta_{k}\right)\left(x_{k}-\lambda_{k} \nabla g\left(x_{k}\right)\right)\right]-P_{C} x^{*}\right\|^{2} \\
& \leq\left\|\left(1-\beta_{k}\right)\left(x_{k}-\lambda_{k} \nabla g\left(x_{k}\right)\right)-x^{*}\right\|^{2} \\
& =\left\|\beta_{k}\left(-x^{*}\right)+\left(1-\beta_{k}\right)\left(x_{k}-\lambda_{k} \nabla g\left(x_{k}\right)-x^{*}\right)\right\|^{2}  \tag{3.12}\\
& \leq \beta_{k}\left\|x^{*}\right\|^{2}+\left(1-\beta_{k}\right)\left\|x_{k}-\lambda_{k} \nabla g\left(x_{k}\right)-x^{*}\right\|^{2} . \tag{3.13}
\end{align*}
$$

By Lemma 2.1(2), we get

$$
\begin{align*}
\left\langle\nabla g\left(x_{k}\right), x_{k}-x^{*}\right\rangle & =\frac{1}{2}\left\langle A^{*}\left(I-P_{Q}\right) A x_{k}+B^{*}\left(I-P_{\tilde{Q}}\right) B x_{k}, x_{k}-x^{*}\right\rangle \\
& =\frac{1}{2}\left[\left\langle A^{*}\left(I-P_{Q}\right) A x_{k}, x_{k}-x^{*}\right\rangle+\left\langle B^{*}\left(I-P_{\tilde{Q}}\right) B x_{k}, x_{k}-x^{*}\right\rangle\right] \\
& =\frac{1}{2}\left[\left\langle\left(I-P_{Q}\right) A x_{k}, A x_{k}-A x^{*}\right\rangle+\left\langle\left(I-P_{\tilde{Q}}\right) B x_{k}, B x_{k}-B x^{*}\right\rangle\right] \\
& \geq \frac{1}{2}\left[\left\|\left(I-P_{Q}\right) A x_{k}\right\|^{2}+\left\|\left(I-P_{\tilde{Q}}\right) B x_{k}\right\|^{2}\right]=2 g\left(x_{k}\right) . \tag{3.14}
\end{align*}
$$

Now using (3.7) and (3.14), we obtain

$$
\begin{align*}
\left\|x_{k}-\lambda_{k} \nabla g\left(x_{k}\right)-x^{*}\right\|^{2} & =\left\|x_{k}-x^{*}\right\|^{2}+\lambda_{k}^{2}\left\|\nabla g\left(x_{k}\right)\right\|^{2}-2 \lambda_{k}\left\langle\nabla g\left(x_{k}\right), x_{k}-x^{*}\right\rangle \\
& \leq\left\|x_{k}-x^{*}\right\|^{2}+\lambda_{k}^{2}\left\|\nabla g\left(x_{k}\right)\right\|^{2}-4 \lambda_{k} g\left(x_{k}\right) \\
& =\left\|x_{k}-x^{*}\right\|^{2}+\frac{\mu_{k}^{2} g^{2}\left(x_{k}\right)}{\left\|\nabla g\left(x_{k}\right)\right\|^{2}}-\frac{4 \mu_{k} g^{2}\left(x_{k}\right)}{\left\|\nabla g\left(x_{k}\right)\right\|^{2}} \\
& =\left\|x_{k}-x^{*}\right\|^{2}-\mu_{k}\left(4-\mu_{k}\right) \frac{g^{2}\left(x_{k}\right)}{\left\|\nabla g\left(x_{k}\right)\right\|^{2}} \tag{3.15}
\end{align*}
$$

Consequently, substituting (3.15) into (3.13) yields

$$
\left\|x_{k+1}-x^{*}\right\|^{2} \leq \beta_{k}\left\|x^{*}\right\|^{2}+\left(1-\beta_{k}\right)\left\|x_{k}-x^{*}\right\|^{2}-\mu_{k}\left(4-\mu_{k}\right)\left(1-\beta_{k}\right) \frac{g^{2}\left(x_{k}\right)}{\left\|\nabla g\left(x_{k}\right)\right\|^{2}}
$$

and (3.10) is obtained. We next show that (3.11) is true. From (3.15), we also have

$$
\begin{align*}
\left\|x_{k}-\lambda_{k} \nabla g\left(x_{k}\right)-x^{*}\right\|^{2} & \leq\left\|x_{k}-x^{*}\right\|^{2}-\mu_{k}\left(4-\mu_{k}\right) \frac{g^{2}\left(x_{k}\right)}{\left\|\nabla g\left(x_{k}\right)\right\|^{2}} \\
& \leq\left\|x_{k}-x^{*}\right\|^{2} . \tag{3.16}
\end{align*}
$$

By using (3.12) and (3.16), we obtain

$$
\begin{aligned}
\left\|x_{k+1}-x^{*}\right\|^{2} \leq & \left\|\beta_{k}\left(-x^{*}\right)+\left(1-\beta_{k}\right)\left(x_{k}-\lambda_{k} \nabla g\left(x_{k}\right)-x^{*}\right)\right\|^{2} \\
= & \beta_{k}^{2}\left\|x^{*}\right\|^{2}+\left(1-\beta_{k}\right)^{2}\left\|x_{k}-\lambda_{k} \nabla g\left(x_{k}\right)-x^{*}\right\|^{2} \\
& +2 \beta_{k}\left(1-\beta_{k}\right)\left\langle x_{k}-\lambda_{k} \nabla g\left(x_{k}\right)-x^{*},-x^{*}\right\rangle \\
\leq & \beta_{k}^{2}\left\|x^{*}\right\|^{2}+\left(1-\beta_{k}\right)^{2}\left\|x_{k}-x^{*}\right\|^{2}+2 \beta_{k}\left(1-\beta_{k}\right)\left\langle x_{k}-x^{*},-x^{*}\right\rangle \\
& +2 \beta_{k}\left(1-\beta_{k}\right) \lambda_{k}\left\langle\nabla g\left(x_{k}\right), x^{*}\right\rangle \\
\leq & \left(1-\beta_{k}\right)\left\|x_{k}-x^{*}\right\|^{2}+\beta_{k}\left[\beta_{k}\left\|x^{*}\right\|^{2}+2\left(1-\beta_{k}\right)\left\langle x_{k}-x^{*},-x^{*}\right\rangle\right. \\
& \left.+2\left(1-\beta_{k}\right) \lambda_{k}\left\langle\nabla g\left(x_{k}\right), x^{*}\right\rangle\right] .
\end{aligned}
$$

This completes the proof.
Lemma 3.6. The sequence $\left\{x_{k}\right\}$ generated by Algorithm 1 is bounded.
Proof. If $\nabla g\left(x_{m}\right)=0$ for some $m \in \mathbb{N}$, then $x_{k}=x_{m}$ for all $k>m$ and hence $\left\{x_{k}\right\}$ is bounded. Assume that $\nabla g\left(x_{k}\right) \neq 0$ for all $k \in \mathbb{N}$. Let $x^{*} \in \Gamma$. Using (3.10), we get

$$
\begin{aligned}
\left\|x_{k+1}-x^{*}\right\|^{2} & \leq \beta_{k}\left\|x^{*}\right\|^{2}+\left(1-\beta_{k}\right)\left\|x_{k}-x^{*}\right\|^{2}-\mu_{k}\left(4-\mu_{k}\right)\left(1-\beta_{k}\right) \frac{g^{2}\left(x_{k}\right)}{\left\|\nabla g\left(x_{k}\right)\right\|^{2}} \\
& \leq \beta_{k}\left\|x^{*}\right\|^{2}+\left(1-\beta_{k}\right)\left\|x_{k}-x^{*}\right\|^{2} \\
& \leq \max \left\{\left\|x^{*}\right\|^{2},\left\|x_{k}-x^{*}\right\|^{2}\right\}
\end{aligned}
$$

By mathematical induction, we deduce that

$$
\left\|x_{k+1}-x^{*}\right\|^{2} \leq \max \left\{\left\|x^{*}\right\|^{2},\left\|x_{1}-x^{*}\right\|^{2}\right\}, \quad \forall k \in \mathbb{N},
$$

it follows that $\left\{x_{k}\right\}$ is bounded.
Now, we are ready to prove a strong convergence theorem of Algorithm 1.
Theorem 3.1. The sequence $\left\{x_{k}\right\}$ generated by Algorithm 1 converges strongly to a solution $x^{*}$ to (1.3) provided that the control sequences $\left\{\beta_{k}\right\}$ and $\left\{\mu_{k}\right\}$ satisfy the following conditions:
(C1) (1) $\lim _{k \rightarrow \infty} \beta_{k}=0$ and (2) $\sum_{k=1}^{\infty} \beta_{k}=\infty$;
(C2) $\inf _{k} \mu_{k}\left(4-\mu_{k}\right)>0$.
Proof. If $\nabla g\left(x_{m}\right)=0$ for some $m \in \mathbb{N}$, then the result is obtained directly by Lemma 3.4. So, we assume that $\nabla g\left(x_{k}\right) \neq 0$ for all $k \in \mathbb{N}$. Let $x^{*}:=P_{\Gamma} 0$. Using (3.10), we get

$$
\left\|x_{k+1}-x^{*}\right\|^{2} \leq \beta_{k}\left\|x^{*}\right\|^{2}+\left(1-\beta_{k}\right)\left\|x_{k}-x^{*}\right\|^{2}-\mu_{k}\left(4-\mu_{k}\right)\left(1-\beta_{k}\right) \frac{g^{2}\left(x_{k}\right)}{\left\|\nabla g\left(x_{k}\right)\right\|^{2}},
$$

it follows that

$$
\begin{equation*}
\mu_{k}\left(4-\mu_{k}\right)\left(1-\beta_{k}\right) \frac{g^{2}\left(x_{k}\right)}{\left\|\nabla g\left(x_{k}\right)\right\|^{2}} \leq \beta_{k}\left\|x^{*}\right\|^{2}+\left\|x_{k}-x^{*}\right\|^{2}-\left\|x_{k+1}-x^{*}\right\|^{2} \tag{3.17}
\end{equation*}
$$

The rest of the proof into will be divided into two cases:
Case 1. Assume that there exists $k_{0} \in \mathbb{N}$ such that $\left\{\left\|x_{k}-x^{*}\right\|\right\}_{k \geq k_{0}}$ is either nonincreasing or nondecreasing. In this case, $\left\{\left\|x_{k}-x^{*}\right\|\right\}$ is convergent because it is bounded. It follows
that $\left\|x_{k}-x^{*}\right\|^{2}-\left\|x_{k+1}-x^{*}\right\|^{2} \rightarrow 0$ as $k \rightarrow \infty$. Then, in view of (3.17) with (C1)(1) and (C2), we obtain

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \frac{g^{2}\left(x_{k}\right)}{\left\|\nabla g\left(x_{k}\right)\right\|^{2}}=0 \tag{3.18}
\end{equation*}
$$

which implies that

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \lambda_{k}\left\|\nabla g\left(x_{k}\right)\right\|=\lim _{k \rightarrow \infty} \frac{g\left(x_{k}\right)}{\left\|\nabla g\left(x_{k}\right)\right\|}=0 \tag{3.19}
\end{equation*}
$$

Let $y_{k}=\left(1-\beta_{k}\right)\left(x_{k}-\lambda_{k} \nabla g\left(x_{k}\right)\right)$. Consider

$$
\left\|x_{k}-y_{k}\right\|=\left\|\left(1-\beta_{k}\right) \lambda_{k} \nabla g\left(x_{k}\right)+\beta_{k} x_{k}\right\| \leq \lambda_{k}\left\|\nabla g\left(x_{k}\right)\right\|+\beta_{k}\left\|x_{k}\right\|,
$$

it follows from (3.19) and (C1)(1) that

$$
\begin{equation*}
\lim _{k \rightarrow \infty}\left\|x_{k}-y_{k}\right\|=0 \tag{3.20}
\end{equation*}
$$

By the same computation as the proof of (3.10), we get

$$
\begin{equation*}
\left\|y_{k}-x^{*}\right\|^{2} \leq \beta_{k}\left\|x^{*}\right\|^{2}+\left(1-\beta_{k}\right)\left\|x_{k}-x^{*}\right\|^{2} . \tag{3.21}
\end{equation*}
$$

For $x \in H_{1}$, we let

$$
h(x):=\left\|\left(I-P_{C}\right) x\right\|^{2}, g^{A}(x):=\frac{1}{2}\left\|\left(I-P_{Q}\right) A x\right\|^{2} \text { and } g^{B}(x):=\frac{1}{2}\left\|\left(I-P_{\tilde{Q}}\right) B x\right\|^{2} .
$$

Using Lemma 2.1(3) and (3.21), we have

$$
\begin{aligned}
h\left(y_{k}\right) & =\left\|y_{k}-P_{C} y_{k}\right\|^{2} \\
& \leq\left\|y_{k}-x^{*}\right\|^{2}-\left\|P_{C} y_{k}-x^{*}\right\|^{2} \\
& \leq \beta_{k}\left\|x^{*}\right\|^{2}+\left\|x_{k}-x^{*}\right\|^{2}-\left\|x_{k+1}-x^{*}\right\|^{2}
\end{aligned}
$$

which implies that

$$
\begin{equation*}
\lim _{k \rightarrow \infty} h\left(y_{k}\right)=0 \tag{3.22}
\end{equation*}
$$

Since $\nabla g^{A}$ and $\nabla g^{B}$ are Lipschitz continuous with coefficients $\|A\|^{2}$ and $\|B\|^{2}$, respectively, one is able to show that $\nabla g$ is Lipschitz continuous with a coefficient $L:=\max \left\{\|A\|^{2},\|B\|^{2}\right\}$. Thus, we have

$$
\left\|\nabla g\left(x_{k}\right)\right\|=\left\|\nabla g\left(x_{k}\right)-\nabla g\left(x^{*}\right)\right\| \leq L\left\|x_{k}-x^{*}\right\|, \quad \forall k \in \mathbb{N}
$$

By the boundedness of $\left\{x_{k}-x^{*}\right\}$, the above inequality yields that $\left\{\nabla g\left(x_{k}\right)\right\}$ is bounded. This together with (3.18) implies that $g\left(x_{k}\right) \rightarrow 0$ as $k \rightarrow \infty$ and hence

$$
\begin{equation*}
\lim _{k \rightarrow \infty} g^{A}\left(x_{k}\right)=\lim _{k \rightarrow \infty} g^{B}\left(x_{k}\right)=0 \tag{3.23}
\end{equation*}
$$

We next show that $\omega_{w}\left(x_{k}\right) \subseteq \Gamma$. Since $\left\{x_{k}\right\}$ is bounded, $\omega_{w}\left(x_{k}\right) \neq \emptyset$. Let $\hat{x} \in \omega_{w}\left(x_{k}\right)$. Then, there exists a subsequence $\left\{x_{k_{j}}\right\}$ of $\left\{x_{k}\right\}$ such that $x_{k_{j}} \rightharpoonup \hat{x}$. Since $g^{A}$ is weakly lower semi-continuous on $H_{1}$, it follows from (3.23) that

$$
0 \leq g^{A}(\hat{x}) \leq \liminf _{j \rightarrow \infty} g^{A}\left(x_{k_{j}}\right)=0
$$

Hence, $g^{A}(\hat{x})=0$, that is, $A \hat{x} \in Q$. Similarly, by using the weakly lower semicontinuity of $g^{B}$ and (3.23), we get $g^{B}(\hat{x})=0$, that is, $B \hat{x} \in \tilde{Q}$. Since $x_{k_{j}} \rightharpoonup \hat{x}$, it also follows from (3.20) that $y_{k_{j}} \rightharpoonup \hat{x}$. By using the weakly lower semicontinuity of $h$ and (3.22), we then deduce that $\hat{x} \in C$. Therefore, $\hat{x} \in \Gamma$ and this means that $\omega_{w}\left(x_{k}\right) \subseteq \Gamma$. Since $x^{*}=P_{\Gamma} 0$, it follows from Lemma 2.1(1) that

$$
\begin{equation*}
\limsup _{k \rightarrow \infty}\left\langle x_{k}-x^{*},-x^{*}\right\rangle=\max _{\hat{x} \in \omega_{w}\left(x_{k}\right)}\left\langle\hat{x}-x^{*},-x^{*}\right\rangle \leq 0 . \tag{3.24}
\end{equation*}
$$

Now using (3.11), we have

$$
\begin{align*}
\left\|x_{k+1}-x^{*}\right\|^{2} \leq & \left(1-\beta_{k}\right)\left\|x_{k}-x^{*}\right\|^{2}+\beta_{k}\left[\beta_{k}\left\|x^{*}\right\|^{2}+2\left(1-\beta_{k}\right)\left\langle x_{k}-x^{*},-x^{*}\right\rangle\right. \\
& \left.+2\left(1-\beta_{k}\right) \lambda_{k}\left\langle\nabla g\left(x_{k}\right), x^{*}\right\rangle\right] \\
\leq & \left(1-\beta_{k}\right)\left\|x_{k}-x^{*}\right\|^{2}+\beta_{k}\left[\beta_{k}\left\|x^{*}\right\|^{2}+2\left(1-\beta_{k}\right)\left\langle x_{k}-x^{*},-x^{*}\right\rangle\right. \\
& \left.+2\left(1-\beta_{k}\right) \lambda_{k}\left\|\nabla g\left(x_{k}\right)\right\|\left\|x^{*}\right\|\right] \\
= & \left(1-\beta_{k}\right)\left\|x_{k}-x^{*}\right\|^{2}+\beta_{k} \delta_{k}, \quad \forall k \in \mathbb{N}, \tag{3.25}
\end{align*}
$$

where

$$
\delta_{k}:=\beta_{k}\left\|x^{*}\right\|^{2}+2\left(1-\beta_{k}\right)\left\langle x_{k}-x^{*},-x^{*}\right\rangle+2\left(1-\beta_{k}\right) \lambda_{k}\left\|\nabla g\left(x_{k}\right)\right\|\left\|x^{*}\right\| .
$$

Using (C1)(1), (3.19), and (3.24), we get $\limsup _{k \rightarrow \infty} \delta_{k} \leq 0$. Recall from (C1)(2) that $\sum_{k=1}^{\infty} \beta_{k}=\infty$. Consequently, by applying Lemma 2.2 to (3.25), we immediately obtain that $x_{k} \rightarrow x^{*}$ as $k \rightarrow \infty$.

Case 2. Assume that $\left\{\left\|x_{k}-x^{*}\right\|\right\}$ is not monotone. There exists a subsequence $\left\{k_{j}\right\}$ of $\{k\}$ such that $\left\|x_{k_{j}}-x^{*}\right\|<\left\|x_{k_{j}+1}-x^{*}\right\|$ for all $j \in \mathbb{N}$. Define a positive interger sequence $\tau(k)$ by

$$
\tau(k):=\max \left\{n \leq k:\left\|x_{n}-x^{*}\right\|<\left\|x_{n+1}-x^{*}\right\|\right\}
$$

for all $k \geq k_{0}$ (for some $k_{0}$ large enough). By Lemma 2.3, $\{\tau(k)\}$ is nondecreasing such that $\tau(k) \rightarrow \infty$ as $k \rightarrow \infty$ and

$$
\begin{equation*}
\left\|x_{\tau(k)}-x^{*}\right\|^{2}-\left\|x_{\tau(k)+1}-x^{*}\right\|^{2} \leq 0 \tag{3.26}
\end{equation*}
$$

for all $k \geq k_{0}$. From (3.17) and (3.26), we have

$$
\begin{aligned}
\mu_{\tau(k)}\left(4-\mu_{\tau(k)}\right)\left(1-\beta_{\tau(k)}\right) \frac{g^{2}\left(x_{\tau(k)}\right)}{\left\|\nabla g\left(x_{\tau(k)}\right)\right\|^{2}} & \leq \beta_{\tau(k)}\left\|x^{*}\right\|^{2}+\left\|x_{\tau(k)}-x^{*}\right\|^{2}-\left\|x_{\tau(k)+1}-x^{*}\right\|^{2} \\
& \leq \beta_{\tau(k)}\left\|x^{*}\right\|^{2}
\end{aligned}
$$

In view of the above inequity with (C1)(1) and (C2), we get

$$
\lim _{k \rightarrow \infty} \frac{g^{2}\left(x_{\tau(k)}\right)}{\left\|\nabla g\left(x_{\tau(k)}\right)\right\|^{2}}=0
$$

By the same way as the proof in Case 1, we obtain

$$
\limsup _{k \rightarrow \infty}\left\langle x_{\tau(k)}-x^{*},-x^{*}\right\rangle=\max _{\hat{x} \in \omega_{w}\left(x_{\tau(k)}\right)}\left\langle\hat{x}-x^{*},-x^{*}\right\rangle \leq 0
$$

and

$$
\begin{equation*}
\left\|x_{\tau(k)+1}-x^{*}\right\|^{2} \leq\left(1-\beta_{\tau(k)}\right)\left\|x_{\tau(k)}-x^{*}\right\|^{2}+\beta_{\tau(k)} \delta_{\tau(k)} \tag{3.27}
\end{equation*}
$$

where
$\delta_{\tau(k)}:=\beta_{\tau(k)}\left\|x^{*}\right\|^{2}+2\left(1-\beta_{\tau(k)}\right)\left\langle x_{\tau(k)}-x^{*},-x^{*}\right\rangle+2\left(1-\beta_{\tau(k)}\right) \lambda_{\tau(k)}\left\|\nabla g\left(x_{\tau(k)}\right)\right\|\left\|x^{*}\right\|$ such that $\limsup _{k \rightarrow \infty} \delta_{\tau(k)} \leq 0$. By looking at (3.27) with the fact that $\left\|x_{\tau(k)}-x^{*}\right\| \leq\left\|x_{\tau(k)+1}-x^{*}\right\|$, we have $\left\|x_{\tau(k)+1}^{k \rightarrow \infty}-x^{*}\right\|^{2} \leq \delta_{\tau(k)}$. This implies that $\limsup _{k \rightarrow \infty}\left\|x_{\tau(k)+1}-x^{*}\right\|^{2} \leq 0$. Consequently, by utilizing Lemma 2.3, we have

$$
0 \leq\left\|x_{k}-x^{*}\right\| \leq\left\|x_{\tau(k)+1}-x^{*}\right\| \rightarrow 0 \text { as } k \rightarrow \infty
$$

Therefore, in both cases we conclude that $\left\{x_{n}\right\}$ converges strongly to $x^{*}=P_{\Gamma} 0$. The proof is complete.

Remark 3.1. It is worth mentioning that there are some advantages of our main result as follows:
(1) If $\left\{x_{k}\right\}$ is a sequence generated by Algorithm 1 such that $\nabla g\left(x_{k}\right) \neq 0$ for all $k \in \mathbb{N}$, then $\left\{x_{k}\right\}$ converges to the minimum-norm solution $x^{*}$ to (1.3), where $x^{*}=P_{\Gamma} 0$.
(2) The choice of stepsize $\lambda_{k}$ defined by (3.7) depends on $x_{k}$ and hence Algorithm 1 does not need to know the value of $\|A\|$ or $\|B\|$.
(3) A result in [35, Theorem 3.1] for solving the SFP is a consequence of Theorem 3.1, namely that if $A=B$ and $Q=\tilde{Q}$ in our problem, then Algorithm 1 is immediately reduced to (1.6) [35, Algorithm 3.1].

Remark 3.2. We note that the concept of choosing the stepsizes $\lambda_{k}$ in (1.5) and (3.7) can be extended to the case of the finite familes of operators $A_{j}$ and sets $Q_{j}(j=1,2, \ldots, n)$ in such a way:

$$
\begin{equation*}
\lambda_{k}=\frac{\mu_{k} g\left(x_{k}\right)}{\left\|\nabla g\left(x_{k}\right)\right\|^{2}} \tag{3.28}
\end{equation*}
$$

where $\mu_{k} \in(0,4)$ and $g: H_{1} \rightarrow \mathbb{R}$ is defined by $g(x):=\frac{1}{2 n} \sum_{j=1}^{n}\left\|\left(I-P_{Q_{j}}\right) A_{j} x\right\|^{2}$ with the gradient given by $\nabla g(x)=\frac{1}{n} \sum_{j=1}^{n} A_{j}^{*}\left(I-P_{Q_{j}}\right) A_{j} x$. It would be interesting to modify the gradient-projection method with the stepsize (3.28) to solve the constrained multiple-set split feasibility problem (CMSSFP) [28] which is formulated as finding a point

$$
x \in \bigcap_{i=1}^{m} C_{i} \text { such that } A_{j} x \in Q_{j}
$$

where $C_{i} \subseteq H_{1}(i=1,2, \ldots, m)$ and $Q_{j} \subseteq H_{2}(j=1,2, \ldots, n)$ are nonempty closed convex subsets and $\left\{A_{j}: H_{1} \rightarrow H_{2}\right\}$ is a finite family of bounded linear operators.

## 4. Self-ADAptive relaxed CQ-TYpe algorithm

Due to our main result in Sect. 3, we consider the two-operator SFP (1.3) for general closed convex subsets $C, Q$, and $\tilde{Q}$; however, finding the explicit forms of the metric projections $P_{C}, P_{Q}$, and $P_{\tilde{Q}}$ in Algorithm 1 may not be easy when these closed convex subsets are complicated. Fortunately, one of the ways for calculating the metric projection onto a sub-level set of a convex function suggested by Fukushima [16] is to compute the sequence of metric projections onto half-spaces containing such a sub-level set. By this idea, Yang [41] considered the SFP (1.1) in the case of two sub-level sets

$$
\begin{equation*}
C=\left\{x \in H_{1}: f_{1}(x) \leq 0\right\} \text { and } Q=\left\{y \in H_{2}: f_{2}(y) \leq 0\right\}, \tag{4.29}
\end{equation*}
$$

where $f_{1}: H_{1} \rightarrow \mathbb{R}$ and $f_{2}: H_{2} \rightarrow \mathbb{R}$ are two convex functions. Also, assume that $f_{1}$ and $f_{2}$ are subdifferentiable on $H_{1}$ and $H_{2}$, respectively, and both $\partial f_{1}$ and $\partial f_{2}$ are bounded operators (i.e., bounded on bounded sets). Yang [41] then introduced the so-called relaxed CQ algorithm for solving the SFP (1.1) constrained by (4.29) as follows:

$$
\left\{\begin{array}{l}
x_{1} \in H_{1},  \tag{4.30}\\
x_{k+1}=P_{C_{k}}\left(x_{k}-\lambda A^{*}\left(I-P_{Q_{k}}\right) A x_{k}\right), \quad k \geq 1,
\end{array}\right.
$$

where $\lambda \in\left(0,2 /\|A\|^{2}\right)$ and $C_{k}$ and $Q_{k}$ are half-spaces given as

$$
C_{k}=\left\{x \in H_{1}: f_{1}\left(x_{k}\right)+\left\langle c_{k}, x-x_{k}\right\rangle \leq 0\right\},
$$

where $c_{k} \in \partial f_{1}\left(x_{k}\right)$ and

$$
Q_{k}=\left\{y \in H_{2}: f_{2}\left(A x_{k}\right)+\left\langle q_{k}, y-A x_{k}\right\rangle \leq 0\right\},
$$

where $q_{k} \in \partial f_{2}\left(A x_{k}\right)$. It follows from the definition of the subdifferential that $C \subseteq C_{k}$ and $Q \subseteq Q_{k}$ for all $k \geq 1$. Since $P_{C_{k}}$ and $P_{Q_{k}}$ have closed forms (see $[6,16]$ ), then the implementation of the relaxed CQ algorithm (4.30) is easier than that of the CQ algorithm (1.2) (in situations that $P_{C}$ and $P_{Q}$ have no closed forms). In addition, López et al. [25, Algorithm 4.1] modified (4.30) by using the self-adaptive stepsize $\lambda_{k}$ (1.5). Vinh et al. [35, Algorithm 4.1] also introduced a relaxation version of the self-adaptive CQ-type algorithm (1.6) to solve this problem.

This section was motivated by the above-mentioned notions and results. We now focus on the two-operator SFP (1.3) in which closed convex subsets $C, Q$, and $\tilde{Q}$ are sub-level sets of convex functions. In what follows, we set the following hypotheses:

- $H_{1}$ and $H_{2}$ are real Hilbert spaces,
- $\emptyset \neq C \subseteq H_{1}$ and $\emptyset \neq Q, \tilde{Q} \subseteq H_{2}$ are given as:

$$
\begin{aligned}
& C=\left\{x \in H_{1}: f_{1}(x) \leq 0\right\}, \\
& Q=\left\{y \in H_{2}: f_{2}(y) \leq 0\right\}, \\
& \tilde{Q}=\left\{y \in H_{2}: \tilde{f}_{2}(y) \leq 0\right\},
\end{aligned}
$$

where $f_{1}: H_{1} \rightarrow \mathbb{R}$ and $f_{2}, \tilde{f}_{2}: H_{2} \rightarrow \mathbb{R}$ are subdifferentiable and convex functions such that their subdifferential operators are bounded,

- $A, B: H_{1} \rightarrow H_{2}$ are two bounded linear operators,
- $\Gamma=\{x \in C: A x \in Q, B x \in \tilde{Q}\} \neq \emptyset$.

Let $x_{k} \in H_{1}$. Denote

$$
\begin{equation*}
C_{k}:=\left\{x \in H_{1}: f_{1}\left(x_{k}\right)+\left\langle c_{k}, x-x_{k}\right\rangle \leq 0\right\}, \tag{4.31}
\end{equation*}
$$

where $c_{k} \in \partial f_{1}\left(x_{k}\right)$,

$$
\begin{equation*}
Q_{k}:=\left\{y \in H_{2}: f_{2}\left(A x_{k}\right)+\left\langle q_{k}, y-A x_{k}\right\rangle \leq 0\right\} \tag{4.32}
\end{equation*}
$$

where $q_{k} \in \partial f_{2}\left(A x_{k}\right)$, and

$$
\begin{equation*}
\tilde{Q}_{k}:=\left\{y \in H_{2}: \tilde{f}_{2}\left(B x_{k}\right)+\left\langle\tilde{q}_{k}, y-B x_{k}\right\rangle \leq 0\right\} \tag{4.33}
\end{equation*}
$$

where $\tilde{q}_{k} \in \partial \tilde{f}_{2}\left(B x_{k}\right)$.
Lemma 4.7. If there exists $x_{k} \in C$ such that $\left\|A^{*}\left(I-P_{Q_{k}}\right) A x_{k}+B^{*}\left(I-P_{\tilde{Q}_{k}}\right) B x_{k}\right\|=0$, then $x_{k} \in \Gamma$.
Proof. Let $x_{k} \in C$ be such that $\left\|A^{*}\left(I-P_{Q_{k}}\right) A x_{k}+B^{*}\left(I-P_{\tilde{Q}_{k}}\right) B x_{k}\right\|=0$. Pick any $p \in \Gamma$. Since $Q \subseteq Q_{k}$ and $\tilde{Q} \subseteq \tilde{Q}_{k}$, then $A p \in Q_{k}$ and $B p \in \tilde{Q}_{k}$. By the same computation as the proof in Lemma 3.4, we get

$$
\begin{aligned}
0 & =\left\|A^{*}\left(I-P_{Q_{k}}\right) A x_{k}+B^{*}\left(I-P_{\tilde{Q}_{k}}\right) B x_{k}\right\|\left\|x_{k}-p\right\| \\
& \geq\left\|\left(I-P_{Q_{k}}\right) A x_{k}\right\|^{2}+\left\|\left(I-P_{\tilde{Q}_{k}}\right) B x_{k}\right\|^{2}
\end{aligned}
$$

which follows that $\left(I-P_{Q_{k}}\right) A x_{k}=\left(I-P_{\tilde{Q}_{k}}\right) B x_{k}=0$ and hence $A x_{k} \in Q_{k}$ and $B x_{k} \in \tilde{Q}_{k}$. By (4.32) and (4.33), we have $f_{2}\left(A x_{k}\right) \leq 0$ and $\tilde{f}_{2}\left(B x_{k}\right) \leq 0$. Thus, $A x_{k} \in Q$ and $B x_{k} \in \tilde{Q}$, i.e., $x_{k} \in \Gamma$.

Using (4.31)-(4.33), a relaxation version of Algorithm 1 is presented as follows.

```
Algorithm 2: Self-adaptive relaxed CQ-type algorithm for the two-operator SFP
    Initialization: Take two real sequences \(\left\{\beta_{k}\right\} \subset(0,1)\) and \(\left\{\mu_{k}\right\} \subset(0,4)\).
    Choose an initial point \(x_{1} \in H_{1}\) arbitrarily and set \(k=1\).
    Iterative Step: Given \(x_{k}\), if \(\left\|A^{*}\left(I-P_{Q_{k}}\right) A x_{k}+B^{*}\left(I-P_{\tilde{Q}_{k}}\right) B x_{k}\right\|=0\), then
    \(x_{k+1}=x_{k}\) and the iterative process stops. Otherwise, calculate
\[
\begin{equation*}
x_{k+1}=P_{C_{k}}\left[\left(1-\beta_{k}\right)\left(x_{k}-\frac{\lambda_{k}}{2}\left(A^{*}\left(I-P_{Q_{k}}\right) A x_{k}+B^{*}\left(I-P_{\tilde{Q}_{k}}\right) B x_{k}\right)\right)\right] . \tag{4.35}
\end{equation*}
\]
Update \(k:=k+1\) and go on to Iterative Step.
```

For the sake of simplicity, we define a function $g_{k}: H_{1} \rightarrow \mathbb{R}$ by

$$
g_{k}(x):=\frac{1}{4}\left(\left\|\left(I-P_{Q_{k}}\right) A x\right\|^{2}+\left\|\left(I-P_{\tilde{Q}_{k}}\right) B x\right\|^{2}\right)
$$

with the gradient given by

$$
\nabla g_{k}(x)=\frac{1}{2}\left(A^{*}\left(I-P_{Q_{k}}\right) A x+B^{*}\left(I-P_{\tilde{Q}_{k}}\right) B x\right), \quad x \in H_{1} .
$$

So, (4.34) and (4.35) become

$$
\lambda_{k}=\frac{\mu_{k} g_{k}\left(x_{k}\right)}{\left\|\nabla g_{k}\left(x_{k}\right)\right\|^{2}} \text { and } x_{k+1}=P_{C_{k}}\left[\left(1-\beta_{k}\right)\left(x_{k}-\lambda_{k} \nabla g_{k}\left(x_{k}\right)\right)\right] .
$$

Below we prove a strong convergence result of Algorithm 2 which extends a result in [35, Theorem 4.1].

Theorem 4.2. Let $\left\{x_{k}\right\}$ be a sequence generated by Algorithm 2 with the control sequences $\left\{\beta_{k}\right\}$ and $\left\{\mu_{k}\right\}$ satisfying:
(C1) (1) $\lim _{k \rightarrow \infty} \beta_{k}=0$ and (2) $\sum_{k=1}^{\infty} \beta_{k}=\infty$;
(C2) $\inf _{k} \mu_{k}\left(4-\mu_{k}\right)>0$.
If $\nabla g_{k}\left(x_{k}\right) \neq 0$ for all $x_{k} \notin C$, then $\left\{x_{k}\right\}$ converges strongly to a point $x^{*} \in \Gamma$.
Proof. If $\nabla g_{m}\left(x_{m}\right)=0$ for some $x_{m} \in C$, then the result is done by Lemma 4.7. So, we suppose that $\nabla g_{k}\left(x_{k}\right) \neq 0$ for all $k \in \mathbb{N}$. Let $x^{*}:=P_{\Gamma} 0$. In view of the proof of Lemma 3.5 with replacing $g$ and $C$ by $g_{k}$ and $C_{k}$, respectively, we deduce that

$$
\begin{equation*}
\left\|x_{k+1}-x^{*}\right\|^{2} \leq \beta_{k}\left\|x^{*}\right\|^{2}+\left(1-\beta_{k}\right)\left\|x_{k}-x^{*}\right\|^{2}-\mu_{k}\left(4-\mu_{k}\right)\left(1-\beta_{k}\right) \frac{g_{k}^{2}\left(x_{k}\right)}{\left\|\nabla g_{k}\left(x_{k}\right)\right\|^{2}} \tag{4.36}
\end{equation*}
$$

and

$$
\begin{align*}
\left\|x_{k+1}-x^{*}\right\|^{2} \leq & \left(1-\beta_{k}\right)\left\|x_{k}-x^{*}\right\|^{2}+\beta_{k}\left[\beta_{k}\left\|x^{*}\right\|^{2}+2\left(1-\beta_{k}\right)\left\langle x_{k}-x^{*},-x^{*}\right\rangle\right. \\
& \left.+2\left(1-\beta_{k}\right) \lambda_{k}\left\langle\nabla g_{k}\left(x_{k}\right), x^{*}\right\rangle\right] . \tag{4.37}
\end{align*}
$$

By (4.36), we obtain that $\left\{x_{k}\right\}$ is bounded and

$$
\begin{equation*}
\mu_{k}\left(4-\mu_{k}\right)\left(1-\beta_{k}\right) \frac{g_{k}^{2}\left(x_{k}\right)}{\left\|\nabla g_{k}\left(x_{k}\right)\right\|^{2}} \leq \beta_{k}\left\|x^{*}\right\|^{2}+\left\|x_{k}-x^{*}\right\|^{2}-\left\|x_{k+1}-x^{*}\right\|^{2} \tag{4.38}
\end{equation*}
$$

Now we consider the rest of the proof into two cases:
Case 1. Suppose that $\left\{\left\|x_{k}-x^{*}\right\|\right\}_{k \geq k_{0}}$ is either nonincreasing or nondecreasing (for some $\left.k_{0}\right)$. We then have $\left\{\left\|x_{k}-x^{*}\right\|\right\}$ is a convergent sequence and so $\left\|x_{k}-x^{*}\right\|^{2}-\left\|x_{k+1}-x^{*}\right\|^{2} \rightarrow$ 0 as $k \rightarrow \infty$. From (4.38), we get

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \frac{g_{k}^{2}\left(x_{k}\right)}{\left\|\nabla g_{k}\left(x_{k}\right)\right\|^{2}}=0 \tag{4.39}
\end{equation*}
$$

which implies that

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \lambda_{k}\left\|\nabla g_{k}\left(x_{k}\right)\right\|=\lim _{k \rightarrow \infty} \frac{g_{k}\left(x_{k}\right)}{\left\|\nabla g_{k}\left(x_{k}\right)\right\|}=0 \tag{4.40}
\end{equation*}
$$

Set $y_{k}=\left(1-\beta_{k}\right)\left(x_{k}-\lambda_{k} \nabla g_{k}\left(x_{k}\right)\right)$. By the same computation as in the proof of Theorem 3.1, we deduce that

$$
\begin{equation*}
\lim _{k \rightarrow \infty}\left\|x_{k}-y_{k}\right\|=0 \tag{4.41}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{k \rightarrow \infty}\left\|\left(I-P_{C_{k}}\right) y_{k}\right\|=0 \tag{4.42}
\end{equation*}
$$

Since $P_{C_{k}} y_{k} \in C_{k}$, it follows from (4.31) and using (4.41), (4.42), and the boundedness assumption on $\partial f_{1}$ that

$$
\begin{align*}
f_{1}\left(x_{k}\right) & \leq\left\langle c_{k}, x_{k}-P_{C_{k}} y_{k}\right\rangle \\
& =\left\langle c_{k}, x_{k}-y_{k}+y_{k}-P_{C_{k}} y_{k}\right\rangle \\
& \leq\left\|c_{k}\right\|\left(\left\|x_{k}-y_{k}\right\|+\left\|\left(I-P_{C_{k}}\right) y_{k}\right\|\right) \rightarrow 0 \text { as } k \rightarrow \infty . \tag{4.43}
\end{align*}
$$

Note that $\nabla g_{k}\left(x^{*}\right)=0$ for all $k \in \mathbb{N}$. Since $\nabla g_{k}$ is Lipschitz continuous with a coefficient $L:=\max \left\{\|A\|^{2},\|B\|^{2}\right\}$, we have

$$
\left\|\nabla g_{k}\left(x_{k}\right)\right\|=\left\|\nabla g_{k}\left(x_{k}\right)-\nabla g_{k}\left(x^{*}\right)\right\| \leq L\left\|x_{k}-x^{*}\right\|, \quad \forall k \in \mathbb{N} .
$$

So, $\left\{\nabla g_{k}\left(x_{k}\right)\right\}$ is bounded. This together with (4.39) yields that $g_{k}\left(x_{k}\right) \rightarrow 0$ as $k \rightarrow \infty$ and hence

$$
\begin{equation*}
\lim _{k \rightarrow \infty}\left\|\left(I-P_{Q_{k}}\right) A x_{k}\right\|=\lim _{k \rightarrow \infty}\left\|\left(I-P_{\tilde{Q}_{k}}\right) B x_{k}\right\|=0 \tag{4.44}
\end{equation*}
$$

Since $P_{Q_{k}}\left(A x_{k}\right) \in Q_{k}$, it follows from (4.32) and using (4.44) and the boundedness assumption on $\partial f_{2}$ that

$$
\begin{align*}
f_{2}\left(A x_{k}\right) & \leq\left\langle q_{k},\left(I-P_{Q_{k}}\right) A x_{k}\right\rangle, \\
& \leq\left\|q_{k}\right\|\left\|\left(I-P_{Q_{k}}\right) A x_{k}\right\| \rightarrow 0 \text { as } k \rightarrow \infty . \tag{4.45}
\end{align*}
$$

Similarly, since $P_{\tilde{Q}_{k}}\left(B x_{k}\right) \in \tilde{Q}_{k}$, it follows from (4.33) and using (4.44) and the boundedness assumption on $\partial \tilde{f}_{2}$ that

$$
\begin{align*}
\tilde{f}_{2}\left(B x_{k}\right) & \leq\left\langle\tilde{q}_{k},\left(I-P_{\tilde{Q}_{k}}\right) B x_{k}\right\rangle, \\
& \leq\left\|\tilde{q}_{k}\right\|\left\|\left(I-P_{\tilde{Q}_{k}}\right) B x_{k}\right\| \rightarrow 0 \text { as } k \rightarrow \infty . \tag{4.46}
\end{align*}
$$

Now let $\hat{x} \in \omega_{w}\left(x_{k}\right)$. Thus, there exists a subsequence $\left\{x_{k_{j}}\right\}$ of $\left\{x_{k}\right\}$ such that $x_{k_{j}} \rightharpoonup \hat{x}$. By the weakly lower semicontinuity of $f_{1}$ and using (4.43), we get

$$
f_{1}(\hat{x}) \leq \liminf _{j \rightarrow \infty} f_{1}\left(x_{k_{j}}\right) \leq 0
$$

This means that $\hat{x} \in C$. Since $A$ and $B$ are bounded linear operators, we also have $A x_{k_{j}} \rightharpoonup$ $A \hat{x}$ and $B x_{k_{j}} \rightharpoonup B \hat{x}$. By the weakly lower semicontinuity of $f_{2}$ and $\tilde{f}_{2}$ and using (4.45), (4.46), we obtain

$$
f_{2}(A \hat{x}) \leq \liminf _{j \rightarrow \infty} f_{2}\left(A x_{k_{j}}\right) \leq 0 \text { and } \tilde{f}_{2}(B \hat{x}) \leq \liminf _{j \rightarrow \infty} \tilde{f}_{2}\left(B x_{k_{j}}\right) \leq 0
$$

which imply that $A \hat{x} \in Q$ and $B \hat{x} \in \tilde{Q}$. Hence, $\hat{x} \in \Gamma$ and so we obtain that $\omega_{w}\left(x_{k}\right) \subseteq \Gamma$. Now, using the characterization of the projection, Lemma 2.1(1) with $P_{\Gamma} 0=x^{*}$, we have

$$
\begin{equation*}
\limsup _{k \rightarrow \infty}\left\langle x_{k}-x^{*},-x^{*}\right\rangle=\max _{\hat{x} \in \omega_{w}\left(x_{k}\right)}\left\langle\hat{x}-x^{*},-x^{*}\right\rangle \leq 0 . \tag{4.47}
\end{equation*}
$$

From (4.37), we get

$$
\begin{aligned}
\left\|x_{k+1}-x^{*}\right\|^{2} \leq & \left(1-\beta_{k}\right)\left\|x_{k}-x^{*}\right\|^{2}+\beta_{k}\left[\beta_{k}\left\|x^{*}\right\|^{2}+2\left(1-\beta_{k}\right)\left\langle x_{k}-x^{*},-x^{*}\right\rangle\right. \\
& \left.+2\left(1-\beta_{k}\right) \lambda_{k}\left\|\nabla g_{k}\left(x_{k}\right)\right\|\left\|x^{*}\right\|\right] \\
= & \left(1-\beta_{k}\right)\left\|x_{k}-x^{*}\right\|^{2}+\beta_{k} \delta_{k}, \quad \forall k \in \mathbb{N},
\end{aligned}
$$

where $\delta_{k}:=\beta_{k}\left\|x^{*}\right\|^{2}+2\left(1-\beta_{k}\right)\left\langle x_{k}-x^{*},-x^{*}\right\rangle+2\left(1-\beta_{k}\right) \lambda_{k}\left\|\nabla g_{k}\left(x_{k}\right)\right\|\left\|x^{*}\right\|$. It follows from (4.40) and (4.47) that $\limsup \delta_{k} \leq 0$. Finally, utilizing Lemma 2.2 with the above inequality, we can conclude that $x_{k} \rightarrow x^{*}$ as $k \rightarrow \infty$.

Case 2. Assume that $\left\{\left\|x_{k}-x^{*}\right\|\right\}$ is not monotone. Using Lemma 2.3 and following the similar argument to the proof in Case 1, one can prove that $\left\{x_{k}\right\}$ also converges strongly to $x^{*}=P_{\Gamma} 0$. So, we omit the proof for this case.

## 5. Numerical experiments

To illustrate the convergence performance of our proposed algorithms and to support our main results, we first employ Algorithm 1 for solving (1.3) in the setting of a Euclidean space (see Example 5.1). After that, we use Algorithm 2 to solve the problem of recovering a sparse signal from a limited number of sampling with two different blurring operations (see Example 5.2). In both examples, we also compare the efficiency of our algorithms with that of some methods based on the operator norms. All the numerical experiments are completed on Apple MacBook Pro with 2 GHz Quad-Core Intel Core i5 with 16 GB memory. The program is implemented in MATLAB R2023a.

Example 5.1. Let $H_{1}=H_{2}=\mathbb{R}^{2}$ with the Euclidean norm. Consider a ball $C$ and a half-space $Q=\tilde{Q}$ given by

$$
C=\left\{(a, b) \in \mathbb{R}^{2}: \sqrt{(a-2)^{2}+b^{2}} \leq 2\right\} \text { and } Q=\left\{(a, b) \in \mathbb{R}^{2}: 3 a+2 b \leq-3\right\}
$$

and two operators $A, B: \mathbb{R}^{2} \rightarrow \mathbb{R}^{2}$ given by $A(a, b)=(-a, 0)$ and $B(a, b)=(0, b)$ for all $(a, b) \in \mathbb{R}^{2}$. One can see that $\Gamma=\{x \in C: A x, B x \in Q\} \neq \emptyset$. We will find the minimumnorm element $x^{*}$ in $\Gamma$ by using our self-adaptive CQ-type algorithm, Algorithm 1. To do this, we arrange the following explicit forms of the metric projections:

$$
P_{C}(a, b)= \begin{cases}(2,0)+\frac{2}{\sqrt{(a-2)^{2}+b^{2}}}(a-2, b), & \text { if }(a, b) \notin C, \\ (a, b), & \text { otherwise }\end{cases}
$$

and

$$
P_{Q}(a, b)= \begin{cases}(a, b)-\frac{3 a+2 b+3}{13}(3,2), & \text { if }(a, b) \notin Q \\ (a, b), & \text { otherwise }\end{cases}
$$

for all $(a, b) \in \mathbb{R}^{2}$. Firstly, we test the convergence behavior of Algorithm 1 by taking $\beta_{k}=\frac{1}{k+1}$ and $\mu_{k}=\frac{2 k}{k+1}$ with the starting point $x_{0}=(4,2)$ as shown in Figure 1. It is observed that $x_{k} \rightarrow(1,-1.5) \in \Gamma$ where $\|(1,-1.5)\|=\min _{p \in \Gamma}\|p\|$.


Figure 1. Illustration of the convergence behavior of Algorithm 1

Next, we analyze the convergence performance of Algorithm 1 by choosing different accelerating sequences $\left\{\mu_{k}\right\}$ and also compare with that of the following algorithms depending on the operator norms.
Algorithm 3: Let $\left\{x_{k}\right\}$ be a sequence generated by (3.8) where

$$
\lambda_{k}:=\lambda \in\left(0, \frac{2}{\max \left\{\|A\|^{2},\|B\|^{2}\right\}}\right) .
$$

Algorithm 4: ([21]) Let $\left\{x_{k}\right\}$ be a sequence generated by (1.4) where $f:=0$.
Each algorithm is equipped with the parameters in Table 1.

TABLE 1. Setting parameters for each algorithm

| Parameters | Algorithm 1 | Algorithm 3 | Algorithm 4 |
| :---: | :---: | :---: | :---: |
| $\beta_{k}=\frac{1}{k+1}$ | $\sqrt{ }$ | $\sqrt{ }$ | $\sqrt{ }$ |
| $\mu_{k}=\frac{\rho k}{k+1}, 0<\rho<4$ | $\sqrt{ }$ | - | - |
| $\delta_{k}=\gamma_{k}=\frac{k}{2 k+2}$ | - | - | $\sqrt{ }$ |
| $0<\lambda<2$ | - | $\sqrt{ }$ | $\sqrt{ }$ |

Table 2. Numerical experiments with the different choices of the stepsizes

|  | Choice of the stepsizes | $k$ (No. of iter.) | CPU time (s) | $x_{k}$ | $E_{k}$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
|  | $\rho=0.5$ | 2198 | 0.2279830 | $(0.9990240,-1.4967059)$ | $9.996 \mathrm{E}-07$ |
| Algorithm 1 | $\rho=1$ | 1100 | 0.1576130 | $(0.9990244,-1.4967074)$ | $9.987 \mathrm{E}-07$ |
|  | $\rho=2$ | 551 | 0.1091000 | $(0.9990253,-1.4967104)$ | $9.968 \mathrm{E}-07$ |
|  | $\rho=3.5$ | 198 | 0.0890400 | $(0.9991338,-1.4968926)$ | $8.726 \mathrm{E}-07$ |
|  | $\rho=3.9$ | 111 | 0.0785620 | $(0.9993079,-1.4966262)$ | $9.585 \mathrm{E}-07$ |
| Algorithm 3 | $\lambda=0.5$ | 5919 | 0.5656190 | $(0.9990239,-1.4967055)$ | $9.998 \mathrm{E}-07$ |
|  | $\lambda=1$ | 2960 | 0.3619410 | $(0.9990240,-1.4967061)$ | $9.995 \mathrm{E}-07$ |
|  | $\lambda=1.9$ | 1558 | 0.2533760 | $(0.9990241,-1.4967063)$ | $9.993 \mathrm{E}-07$ |
| Algorithm 4 ([21]) | $\lambda=0.5$ | 11837 | 1.1435740 | $(0.9990238,-1.4967052)$ | $9.999 \mathrm{E}-07$ |
|  | $\lambda=1$ | 5919 | 0.6078450 | $(0.9990239,-1.4967055)$ | $9.998 \mathrm{E}-07$ |
|  | $\lambda=1.9$ | 3115 | 0.3736660 | $(0.9990238,-1.4967052)$ | $9.998 \mathrm{E}-07$ |

We choose the starting point $x_{0}=x_{1}=(2,2)$ and use the stopping criterion for the iterative process as: $E_{k}:=g\left(x_{k}\right)<10^{-6}$, where $g$ is defined by (3.9). Now the comparision of the numerical experiments of Algorithms 1,3, and 4 are shown in Table 2.
Remark 5.3. By testing the performance of Algorithms 1,3, and 4 and from Table 2, we observe that
(1) All studied algorithms give the approximate solutions close to $(1,-1.5)$ which is the minimum-norm solution.
(2) Algorithm 1 converges the fastest and takes the least time.
(3) The choice of the stepsizes influencses the convergence behavior of all studied algorithms. Namely that if $\left\{\mu_{k}\right\}$ is taken close to 4 (for Algorithm 1) and $\lambda$ is taken close to 2 (for Algorithms 3 and 4), then the number of iterations and the CPU time have reduction. Meanwhile, choosing different starting points has no significant impact on their convergence behavior.

Example 5.2. (Compressed Sensing $[25,30])$. Here, we consider the problem of recovering a sparse signal $x \in \mathbb{R}^{N}$ from the observation of two signals $y, \tilde{y} \in \mathbb{R}^{M}(M<N)$ via the linear equation systems:

$$
\begin{equation*}
y=A x+\varepsilon \text { and } \tilde{y}=B x+\tilde{\varepsilon}, \tag{5.48}
\end{equation*}
$$

where $A, B: \mathbb{R}^{N} \rightarrow \mathbb{R}^{M}$ are two bounded linear observation operators (they are often ill-condition) and $\varepsilon, \tilde{\varepsilon}$ are additive noises. The problem (5.48) can be solved by using the LASSO technique ([34]) in the froms of the constrained least-squares problem:

$$
\begin{equation*}
\operatorname{minimize} \frac{1}{2}\|A x-y\|_{2}^{2} \text { and } \frac{1}{2}\|B x-\tilde{y}\|_{2}^{2} \tag{5.49}
\end{equation*}
$$

with respect to $x \in C:=\left\{x \in \mathbb{R}^{N}:\|x\|_{1} \leq t\right\}$, where $t>0$ is a given constant. If (5.49) has a solution, we see that (5.49) is a particular case of the two-operator SFP (1.3) where $Q=\{y\}$ and $\tilde{Q}=\{\tilde{y}\}$. Since $C$ is the closed $l_{1}$ ball in $\mathbb{R}^{N}$ with the radius $t$, we will employ
the relaxation version of our self-adaptive CQ-type algorithm, Algorithm 2 to solve (5.49). Define $f_{1}(x)=\|x\|_{1}-t$ and consider the half-space $C_{k}$ denoted by (4.31). The closed form of the metric projection from $\mathbb{R}^{N}$ onto $C_{k}$ is as follows:

$$
P_{C_{k}}(x)= \begin{cases}x, & \text { if } f_{1}\left(x_{k}\right)+\left\langle c_{k}, x-x_{k}\right\rangle \leq 0 \\ x-\frac{f_{1}\left(x_{k}\right)+\left\langle c_{k}, x-x_{k}\right\rangle}{\left\|c_{k}\right\|^{2}} c_{k}, & \text { otherwise }\end{cases}
$$

where $c_{k} \in \partial f_{1}\left(x_{k}\right)$ is chosen as

$$
c_{k}^{(i)}= \begin{cases}1, & \text { if } x_{k}^{(i)}>0 \\ 0, & \text { if } x_{k}^{(i)}=0 \\ -1, & \text { if } x_{k}^{(i)}<0\end{cases}
$$

see [17, Section 5].
In our experiment, two sampling matrices $A, B \in \mathbb{R}^{M \times N}$ are generated randomly from normal distributions with $N=2048$ and $M=1024$. The sparse signal $x^{*} \in \mathbb{R}^{N}$ is generated from a uniform distribution in $[-2,2]$ with $m$ nonzero components. The measured values $y$ and $\tilde{y}$ are generated by white Gaussian noise with the signal-to-noise ratio (SNR) as 40 and 50 decibles, respectively. Set $t=m$. We test three cases as follows:

$$
\text { Case 1: } m=10, \quad \text { Case 2: } m=50, \quad \text { Case 3: } m=100
$$

We compare the signal recovery performance of Algorithm 2 with that of the following algorithm depending on $\|A\|$ and $\|B\|$.

Algorithm 5: Let $\left\{x_{k}\right\}$ be a sequence generated by (4.35) where

$$
\lambda_{k}:=\lambda \in\left(0, \frac{2}{\max \left\{\|A\|^{2},\|B\|^{2}\right\}}\right) .
$$

Let $\beta_{k}=\frac{1}{k+1}$ and $\mu_{k}=\frac{2 k}{k+1}$ for Algorithm 2 and $\beta_{k}=\frac{1}{k+1}$ and $\lambda=\frac{1}{\max \left\{\|A\|^{2},\|B\|^{2}\right\}}$ for Algorithm 5. The process is started with the initial signal $x_{1}=0$. The restoration accuracy is measured by the mean squared error (MSE), i.e.,

$$
\operatorname{MSE}(k)=\frac{1}{N}\left\|x^{*}-x_{k}\right\|^{2}<10^{-4}
$$

where $x^{*}$ is the original signal and $x_{k}$ is an estimated signal of $x^{*}$. Now, the numerical results of recovering the signal $x^{*}$ are reported as Figures 2-7.

Remark 5.4. By the simple experiments as shown in Figures 2-7, we note that
(1) The original signals $x^{*}$ can be recovered by Algorithms 2 and 5.
(2) If the number of spikes of $x^{*}$ increases, then both methods also require an increase in the number of iterations and the CPU time. However, the number of iterations and the CPU time of using Algorithm 2 are less than those of using Algorithm 5.


Figure 2. Signal recovery experiment in Case 1.
From top to bottom: original signal; observation data using $A$; observation data using $B$; recovered signal by Algorithm 2; recovered signal by Algorithm 5


Figure 3. The mean squared error versus the number of iterations in Case 1


Figure 4. Signal recovery experiment in Case 2.
From top to bottom: original signal; observation data using $A$; observation data using $B$; recovered signal by Algorithm 2; recovered signal by Algorithm 5


Figure 5. The mean squared error versus the number of iterations in Case 2


Figure 6. Signal recovery experiment in Case 3.
From top to bottom: original signal; observation data using $A$; observation data using $B$; recovered signal by Algorithm 2; recovered signal by Algorithm 5


Figure 7. The mean squared error versus the number of iterations in Case 3

For the closed forms of some metric projections onto simple closed convex subsets in Hilbert spaces, the reader is referred to [6, Chapter 4]. There are also some examples for the split feasibility problem and related problems in the infinite-dimensional Hilbert spaces, see [23, 33, 35].

## CONCLUSION

This paper discusses and analyzes the convergence results on the two-operator split feasibility problem (two-operator SFP) in Hilbert spaces, namely finding a point of a closed convex subset of a Hilbert space such that each of its images under two given bounded linear operators belongs to a closed convex subset of another Hilbert space. We introduce a self-adaptive CQ-type algorithm where the stepsize does not depend on such bounded linear operator norms. Under some mild conditions, we then prove that the sequence generated by the proposed algorithm converges strongly to the minimum-norm solution of the two-operator SFP. A relaxation version of our proposed algorithm is also introduced for solving the problem constrained by sub-level sets of convex functions. Our main results improve the result of Kangtunyakarn [21, Theorems 3.1] in terms of selecting the stepsize in the algorithm and generalize the results of Vinh et al. [35, Theorems 3.1 and 4.1] for the split feasibility problem (also improve the results of Xu [40], Wang and Xu [37], Yao et al. [43] and Chuang [7]). In addition, it is observed from our numerical experiments that our self-adaptive CQ-type algorithms (without any operator norms) are more efficient than the CQ-type algorithms based on the operator norms.

Acknowledgments. W. Khuangsatung would like to thank Rajamangala University of Technology Thanyaburi (RMUTT) under Thailand Science Research and Innovation Promotion Funding (TSRI) (Contract No. FRB660012/0168 and under project number FRB66E0635) for financial support. P. Jailoka would like to thank Thailand Science Research and Innovation Fund and University of Phayao (Grant FF67-UoE).

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[^0]:    Received: 30.08.2022. In revised form: 17.07.2023. Accepted: 24.07.2023
    2010 Mathematics Subject Classification. 47H10, 90C25.
    Key words and phrases. split feasibility problems, CQ algorithms, self-adaptive stepsizes, strong convergence.
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