# Improving the speed of convergence in the method of projections onto convex sets 

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#### Abstract

A serious drawback of the method of projections onto convex sets to find a point in the intersection of a finite number of closed convex sets in an Euclidean space, is its often very slow convergence. This bad behaviour, sometimes called the "tunneling effect", seems a.o. to be connected with the monotone behaviour of the usual algorithms. We present a new algorithm that can interrupt at different steps this monotone behaviour; this can substantially improve the speed of convergence.


## 1. Introduction

Finding a point in the intersection of a finite number of closed convex sets $\left\{C_{j}\right\}_{j=1}^{r}$ in an $m$-dimensional Euclidean space $\mathbb{R}^{m}(\mathbb{R}$ the set of reals), with $C^{*} \equiv \bigcap_{j=1}^{r} C_{j} \neq \emptyset$, is a problem that often arises in applied mathematics; we refer to [7] and [9] for a short overview of some general applications, and to [15] and [16] for specific applications in image processing. The stated problem is also known as "solving the convex feasibility problem".

The method of projections onto convex sets (abbreviated as POCS), seems to be very well suited to solve this problem. Originating from the work by J. von Neumann [14] for alternating orthogonal projections onto closed linear subspaces, it has been generalized and extended to cases involving closed convex subsets and to parallel methods for the algorithm; e.g., see [1], [3], [4] and the references in these papers. Acceleration schemes for POCS-like methods for finding the shortest-distance projection of a

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given point onto the intersection of a number of closed linear varieties have been given in [8]; in [10], an acceleration scheme for Kaczmarz's method is presented.

The POCS method, however, has one serious drawback: it often converges very slowly. This fact, sometimes called the "tunneling effect" [11], is in an intuitive manner explained as follows: in order to reach the intersection, the constructed iteration sequence enters a narrow corridor (a "tunnel") formed by two of the involved sets, and at each iteration step the progress towards the intersection set, caused by projecting the current iteration point onto the involved sets, is very small.

This explanation, however, is only half of the truth. Indeed, the decisive factor concerning the speed of convergence seems to be the combination of the structure of the involved sets together with the starting point of the iteration sequence (this fact makes it difficult to compare the rate of convergence of different algorithms in practical applications on the base of theoretical results). By way of example, we consider the problem of finding, with the POCS method, a point in the intersection of the following twelve disks $\left\{C_{j}\right\}_{j=1}^{12}$

$$
\begin{equation*}
\left(x-\cos \left(\frac{j \pi}{12}\right)\right)^{2}+\left(y-\sin \left(\frac{j \pi}{12}\right)\right)^{2} \leq 1, \quad \text { for } j=1,2, \ldots, 12 \tag{1}
\end{equation*}
$$

in the plane (clearly, $(0,0)$ is a point in the intersection).
Denoting by $P_{j}$ the projection onto the corresponding disk $C_{j}$, denoting momentarily the points in $\mathbb{R}^{m}$ by boldface letters, the iteration sequence as $\left\{\boldsymbol{x}_{k}\right\}_{k=0}^{+\infty}$ with $\boldsymbol{x}_{k} \equiv\left(x_{k}, y_{k}\right)$, and using the iteration $\boldsymbol{x}_{k+1}=T \boldsymbol{x}_{k}$ with $T \equiv P_{12} P_{11} \ldots P_{2} P_{1}$, we obtain very different convergence results depending on the choice of the starting point $\boldsymbol{x}_{\boldsymbol{o}} \equiv\left(x_{0}, y_{0}\right)$ in the plane. For $\left(x_{0}, y_{0}\right) \equiv(-3,0)$, a point in the intersection is obtained after one iteration; on the other hand, for $\left(x_{0}, y_{0}\right) \equiv(3,4)$, the sum of the distances of the current iteration point to the sets $C_{j}$ is $3.661634 \times 10^{-3}, 5.49556 \times 10^{-4}$ and $1.66893 \times 10^{-5}$ after 25,50 and 100 iterations respectively. A behaviour comparable to this one is observed when using different existing adaptations of the method of pure projections, as for instance for the parallel method

$$
\begin{equation*}
\boldsymbol{x}_{k+1}=\boldsymbol{x}_{k}+\lambda_{k+1} \sum_{j=1}^{r} \mu_{k+1}(j)\left(P_{j} \boldsymbol{x}_{k}-\boldsymbol{x}_{k}\right), \tag{2}
\end{equation*}
$$

with one variable relaxation coefficient $\lambda_{k+1}$, and with nonnegative weights $\mu_{k+1}(j)(j=1, \ldots, r)$ such that $\sum_{j=1}^{r} \mu_{k+1}(j)=1$.

Due to the fact that, when convex sets $C_{j}$ are given, there is usually no clear indication of how to choose the starting point of the iteration in order to avoid the tunneling effect, it seems that, in order to circumvent this effect, it might be favourable to interrupt at different iteration steps during the procedure the "monotone" way of converging. By this we mean that, contrary to the existing algorithms where it is true that for the iteration sequence $\left\{\boldsymbol{x}_{n}\right\}_{n=0}^{+\infty}$ in $\mathbb{R}^{m}$ we have that

$$
\begin{equation*}
\left\|\boldsymbol{x}_{k+1}-\boldsymbol{w}\right\| \leq\left\|\boldsymbol{x}_{k}-\boldsymbol{w}\right\|, \forall \boldsymbol{w} \in C^{*}, \quad \text { for all } k \tag{3}
\end{equation*}
$$

the speed of convergence may be improved by an algorithm that automatically chooses, at some instances, the update iteration point $\boldsymbol{x}_{k+1}$ such that (3) is not necessarily true. At first glance, this may seem rather awkward, but in this way we may approach the tunnel from another direction, leading to an improvement of the convergence speed in those cases where the tunneling effect is strong, while keeping an acceptable speed of convergence in the other cases. With the algorithm presented in this paper, convergence to a point in the intersection of the given example is obtained after 9 iterations when starting from $\left(x_{0}, y_{0}\right) \equiv(-3,0)$, and after 9 iterations when starting from $\left(x_{0}, y_{0}\right) \equiv(3,4)$.

The method we present is an adaptation of the parallel method (2) with fixed weights $\mu_{k+1}(j)$. As has been shown by Pierra [12], the parallel method (2) in $\mathbb{R}^{m}$ may then be seen as a sequential method in the product space $\left(\mathbb{R}^{m}\right)^{r}$ (extensions of this with variable weights have been given in [5]). According to the Pierra method, in $\left(\mathbb{R}^{m}\right)^{r}$ there are two important subsets: the linear subspace $\mathcal{D}$, which is the canonical imbedding $q\left(\mathbb{R}^{m}\right)$ of $\mathbb{R}^{m}$ into $\left(\mathbb{R}^{m}\right)^{r}$, and the closed convex set $\mathcal{F} \equiv C_{1} \times C_{2} \times \cdots \times C_{r}$. Denoting for the moment the points in $\left(\mathbb{R}^{m}\right)^{r}$ by capital letters (in order to distinguish them from the points in $\mathbb{R}^{m}$ ), the ordinary parallel projection method (2) in $\mathbb{R}^{m}$ with fixed equal weights $\mu_{k+1}(j) \equiv \frac{1}{r}$ at each iteration step is then equivalent with the following procedure in $\left(\mathbb{R}^{m}\right)^{r}$ starting from some point $X_{0} \equiv q\left(\boldsymbol{x}_{0}\right)$ in $\mathcal{D}$ : whenever $X_{k}$ in $\mathcal{D}$ has been obtained, the next iteration point $X_{k+1}$ in $\mathcal{D}$ is obtained as a result of the following two steps: put

$$
\begin{equation*}
Y_{k+1}=X_{k}+\lambda_{k+1}\left(P_{\mathcal{F}} X_{k}-X_{k}\right), \tag{4}
\end{equation*}
$$

where $P_{\mathcal{F}} X_{k}$ is the projection of $X_{k}$ onto $\mathcal{F}$; then put

$$
\begin{equation*}
X_{k+1}=P_{\mathcal{D}}\left(Y_{k+1}\right)=X_{k}+\lambda_{k+1}\left(P_{\mathcal{D}}\left(P_{\mathcal{F}} X_{k}\right)-X_{k}\right) . \tag{5}
\end{equation*}
$$

Then, for suitable $\lambda_{k+1}$ the sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$ in $\mathcal{D}$ is convergent to a point of $\mathcal{F} \cap \mathcal{D}$, and for each point $W$ in $\mathcal{F} \cap \mathcal{D}$ the Fejér-monotone property ([13]) analogous to (3) is fulfilled (e.g., see [6]).

Formula (5), valid in $\mathcal{D} \subset\left(\mathbb{R}^{m}\right)^{r}$, has the typical structure of an iteration process: the update $X_{k+1}$ of $X_{k}$ is found by starting from $X_{k}$ and choosing a direction in $\mathcal{D}$ (determined by $P_{\mathcal{D}}\left(P_{\mathcal{F}} X_{k}\right)$ ) and a steplength $\lambda_{k+1}$. This update process in (5) is derived from the update given in (4) which, however, is a process in $\left(\mathbb{R}^{m}\right)^{r}$ and as such gives much more flexibility. Indeed, while in (4) the computation of $Y_{k+1}$ is started from $X_{k} \in \mathcal{D}$ in order to obtain (5), there is certainly no need to restrict such starting point of (4) to $\mathcal{D}$. When taking as starting point for (4) a point $Z_{k} \in\left(\mathbb{R}^{m}\right)^{r}$ that does not necessarily belong to $\mathcal{D}$ but is such that $P_{\mathcal{D}} Z_{k}=$ $X_{k}$, then formulas (4) and (5) are changed into

$$
\begin{align*}
Y_{k+1} & =Z_{k}+\lambda_{k+1}\left(P_{\mathcal{F}} Z_{k}-Z_{k}\right)  \tag{4'}\\
X_{k+1} & =P_{\mathcal{D}}\left(Y_{k+1}\right)=X_{k}+\lambda_{k+1}\left(P_{\mathcal{D}}\left(P_{\mathcal{F}} Z_{k}\right)-X_{k}\right) \tag{5’}
\end{align*}
$$

and so also ( $5^{\prime}$ ) is a typical iteration process in $\mathcal{D}$, for a suitable step-length $\lambda_{k+1}$ and a direction determined by $P_{\mathcal{D}}\left(P_{\mathcal{F}} Z_{k}\right)$.

The basic idea in our paper is precisely to use the freedom of choosing $Z_{k}$ in (4'), which itself entails more possibilities for the value of $\lambda_{k+1}$. Just for illustrative purposes we sketch the first two steps of our iteration algorithm. As concerns notation, when $Z_{k} \in\left(\mathbb{R}^{m}\right)^{r}$ but $Z_{k} \notin \mathcal{F}$, we denote by $\mathcal{P}_{k+1}$ the hyperplane supporting $\mathcal{F}$ at $P_{\mathcal{F}} Z_{k}$.

Let $Z_{0} \equiv X_{0}$ be a starting point in $\mathcal{D}$. Put

$$
\begin{aligned}
Y_{1} & =Z_{0}+\lambda_{1}\left(P_{\mathcal{F}} Z_{0}-Z_{0}\right) \\
X_{1} & =P_{\mathcal{D}} Y_{1}
\end{aligned}
$$

where $\lambda_{1}$ is determined such that $X_{1} \in \mathcal{P}_{1} \cap \mathcal{D}$ (such $\lambda_{1}$ always exists). This is the first iteration point of our sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$ in $\mathcal{D}$. In order to obtain $X_{2}$, we construct $Z_{1}$ by using a relaxed projection of $Y_{1}$ onto $\mathcal{D}$ by means of

$$
\begin{equation*}
Z_{1}=Y_{1}+\left(1+\gamma_{1}\right)\left(X_{1}-Y_{1}\right)=X_{1}+\gamma_{1}\left(X_{1}-Y_{1}\right) \tag{6}
\end{equation*}
$$

for some $\gamma_{1}>0$.

For this $Z_{1}$ we can verify if there is a real number $\lambda_{2}$ such that the projection onto $\mathcal{D}$ of the point $Z_{1}+\lambda_{2}\left(P_{\mathcal{F}} Z_{1}-Z_{1}\right)$ belongs to $\mathcal{P}_{2} \cap \mathcal{D}$. If $\lambda_{2}$ exists and is bigger than 1 , we put

$$
\begin{align*}
Y_{2} & =Z_{1}+\lambda_{2}\left(P_{\mathcal{F}} Z_{1}-Z_{1}\right),  \tag{7}\\
X_{2} & =P_{\mathcal{D}} Y_{2} \tag{8}
\end{align*}
$$

(which is the second iteration point in this case),

$$
\begin{equation*}
Z_{2}=Y_{2}+\left(1+\gamma_{2}\right)\left(X_{2}-Y_{2}\right) \quad \text { for some } \gamma_{2}>0 \tag{9}
\end{equation*}
$$

(which is the starting point for the next projection onto $\mathcal{F}$ in this case).
On the other hand, if $\lambda_{2} \leq 1$, we put

$$
\begin{align*}
& Y_{2}=P_{\mathcal{F}} Z_{1},  \tag{10}\\
& Z_{2} \equiv X_{2}=P_{\mathcal{D}} Y_{2} \tag{11}
\end{align*}
$$

(which is now at the same time the second iteration point $X_{2}$ of the approximation sequence, and the starting point for the next projection onto $\mathcal{F}$ ).

Finally, if $\lambda_{2}$ does not exist, then $X_{1}$ is chosen as the starting point to obtain $X_{2}$, in the same way as $X_{1}$ was constructed from $X_{0}$.

So, from $X_{1}$ on, each next iteration point $X_{k+1}$ of the sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$ can be obtained on the discriminating base of the value of some parameter $\lambda_{k+1}$, and this may cause a different behaviour with respect to Fejér-monotony; in particular, the (monotone but) slow convergence may be interrupted at some iteration points.

The values of $\lambda_{k+1}$ involved in this procedure play the role of relaxation coefficients connected with projections onto $\mathcal{F}$ when they are bigger than 1 ; as they may also be bigger than 2 , the corresponding relaxed projection onto $\mathcal{F}$ has no longer the non-expansivity property.

The rest of the paper is organized as follows. In Section 2 we introduce the results that are needed to show that the algorithm is meaningful, and we describe the algorithm. In Section 3 we show that the sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$ in $\mathcal{D} \subset\left(\mathbb{R}^{m}\right)^{r}$, constructed according to the algorithm, is convergent to a point of $\mathcal{F} \cap \mathcal{D}$, leading to convergence of the corresponding sequence $\left\{\boldsymbol{x}_{k}\right\}_{k=0}^{+\infty}$ in $\mathbb{R}^{m}$ to a point in $C^{*}$. Finally, in Section 4 we illustrate the new algorithm for some examples, and we compare its behaviour for different starting points to the method of pure projections and to the parallel method (2).

## 2. Description of the algorithm

Let $\mathbb{R}^{m}$ be the $m$-dimensional Euclidean space with standard inner product $\langle$,$\rangle and norm \mid$ derived from $\langle$,$\rangle ; denote \left(\mathbb{R}^{m},\langle\rangle,,| |\right)$ for short by $H$. Elements of $H$ are denoted by boldface letters.

Suppose that in $H, r$ closed convex subsets $\left\{C_{j}\right\}_{j=1}^{r}$ are given, having nonempty intersection $C^{*} \equiv \bigcap_{j=1}^{r} C_{j} \neq \emptyset$. Projection onto $C_{j}$ is denoted as $P_{j}$. We want to obtain, in an iterative manner, a point in $C^{*}$.

Consider the $r$-fold product $\left(\mathbb{R}^{m}\right)^{r}$ of $\mathbb{R}^{m}$; elements of $\left(\mathbb{R}^{m}\right)^{r}$ are denoted by capital letters. We introduce an inner product $\langle\langle$,$\rangle and$ norm \| \| as follows: when $V \equiv\left(\boldsymbol{v}_{1}, \boldsymbol{v}_{2}, \ldots \boldsymbol{v}_{r}\right) \in\left(\mathbb{R}^{m}\right)^{r}$ and $W \equiv$ $\left(\boldsymbol{w}_{1}, \boldsymbol{w}_{2}, \ldots, \boldsymbol{w}_{r}\right) \in\left(\mathbb{R}^{m}\right)^{r}$, put $\langle\langle V, W\rangle\rangle=\sum_{j=1}^{r}\left\langle\boldsymbol{v}_{j}, \boldsymbol{w}_{j}\right\rangle,\|V\|^{2}=\sum_{j=1}^{r}\left|\boldsymbol{v}_{j}\right|^{2}$.

We denote $\left(\left(\mathbb{R}^{m}\right)^{r},\langle\langle\rangle\rangle,,\| \|\right)$ for short by $\mathcal{H}$.
In $\mathcal{H}$ we consider the subsets $\mathcal{D}$ and $\mathcal{F}$, defined as follows. $\mathcal{D}$ is the set of all $r$-tuples with equal components, i.e., for $\boldsymbol{v} \in H$ we have that $(\boldsymbol{v}, \boldsymbol{v}, \ldots, \boldsymbol{v}) \in \mathcal{D} \subset \mathcal{H} . \mathcal{D}$ is the image of $H$ under the canonical imbedding $q: H \rightarrow \mathcal{H}$, where for $\boldsymbol{v} \in H$ we put $q(\boldsymbol{v}) \equiv(\boldsymbol{v}, \boldsymbol{v}, \ldots, \boldsymbol{v})$. $\mathcal{D}$ is a closed linear subspace of $\mathcal{H}$. Projection onto $D$ is denoted as $P_{\mathcal{D}}$.

The subset $\mathcal{F}$ of $\mathcal{H}$ is defined as the Cartesian product of the convex sets $\left\{C_{j}\right\}_{j=1}^{r}$ in $H$, i.e., $\mathcal{F}=C_{1} \times C_{2} \times \cdots \times C_{r}$. It is a closed convex subset of $\mathcal{H}$. Projection onto $\mathcal{F}$ is denoted as $P_{\mathcal{F}}$. Clearly, as $C^{*} \neq \emptyset$ we also have that $\mathcal{F} \cap \mathcal{D} \neq \emptyset$, and, moreover, $q\left(C^{*}\right)=\mathcal{F} \cap \mathcal{D}$. Hence, obtaining a point in $C^{*} \subset H$ is equivalent to obtaining a point in $\mathcal{F} \cap \mathcal{D} \subset \mathcal{H}$. So, our aim is to construct a sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$ in $\mathcal{D} \subset \mathcal{H}$, starting from some given point $X_{0} \in \mathcal{D}$, that converges in $\mathcal{H}$ to a point in $\mathcal{F} \cap \mathcal{D}$; the corresponding sequence $\left\{\boldsymbol{x}_{k}\right\}_{k=0}^{+\infty}$ in $H$ with $q\left(\boldsymbol{x}_{k}\right)=X_{k}$ will then be convergent to a point in $C^{*}$.

In order to construct the sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$ in $\mathcal{H}$, we need some preliminary results.

Lemma 1. Let $V \equiv\left(\boldsymbol{v}_{1}, \boldsymbol{v}_{2}, \ldots, \boldsymbol{v}_{r}\right) \in \mathcal{H}$. Then
(i) $P_{\mathcal{F}} V=\left(P_{1} \boldsymbol{v}_{1}, P_{2} \boldsymbol{v}_{2}, \ldots, P_{r} \boldsymbol{v}_{r}\right)$
(ii) $P_{\mathcal{D}} V=q\left(\frac{1}{r} \sum_{k=1}^{r} \boldsymbol{v}_{k}\right)$

Proof. (i) For arbitrary $S \equiv\left(s_{1}, s_{2}, \ldots, s_{r}\right) \in \mathcal{F}$ we have

$$
\|V-S\|^{2}=\sum_{j=1}^{r}\left|\boldsymbol{v}_{j}-\boldsymbol{s}_{j}\right|^{2}
$$

and this is minimal when $\boldsymbol{s}_{j}=P_{j} \boldsymbol{v}_{j}$.
(ii) See [12, Lemma 1.1].

For a given point $Z$ in $\mathcal{H}$ with projection $P_{\mathcal{F}} Z$ onto $\mathcal{F}$ (assume $Z \neq$ $P_{\mathcal{F}} Z$ ), consider the hyperplane $\mathcal{P}$ going through $P_{\mathcal{F}} Z$ which is orthogonal onto $Z-P_{\mathcal{F}} Z$, i.e.,

$$
\begin{equation*}
\mathcal{P}=\left\{X \in \mathcal{H}:\left\langle\left\langle X-P_{\mathcal{F}} Z, P_{\mathcal{F}} Z-Z\right\rangle\right\rangle=0\right\} . \tag{12}
\end{equation*}
$$

On the straight line through $Z$ and $P_{\mathcal{F}} Z$ we consider the points $Y$ given by the following expression

$$
\begin{equation*}
Y=Z+\lambda\left(P_{\mathcal{F}} Z-Z\right), \tag{13}
\end{equation*}
$$

where $\lambda$ is a real parameter.
Lemma 2. Let $Z \in \mathcal{H}, Z \neq P_{\mathcal{F}} Z$, and put $Y=Z+\lambda\left(P_{\mathcal{F}} Z-Z\right)$. Then $P_{\mathcal{D}} Y$ belongs to $\mathcal{P}$ if and only if $\lambda$ is given by

$$
\begin{equation*}
\lambda=\frac{\left\langle\left\langle P_{\mathcal{D}} Z-P_{\mathcal{F}} Z, Z-P_{\mathcal{F}} Z\right\rangle\right\rangle}{\left\|P_{\mathcal{D}}\left(P_{\mathcal{F}} Z\right)-P_{\mathcal{D}} Z\right\|^{2}} . \tag{14}
\end{equation*}
$$

Proof. Assume that there exists a real number $\lambda$ such that $P_{\mathcal{D}} Y$ belongs to $\mathcal{P}$.

As $P_{\mathcal{D}}$ is linear, $P_{\mathcal{D}} Y$ is given by

$$
\begin{equation*}
P_{\mathcal{D}} Y=P_{\mathcal{D}} Z+\lambda\left(P_{\mathcal{D}}\left(P_{\mathcal{F}} Z\right)-P_{\mathcal{D}} Z\right) \tag{15}
\end{equation*}
$$

In the expression (12) of $\mathcal{P}$ we write the first part in the inner product sign as $X-P_{\mathcal{D}} Z+P_{\mathcal{D}} Z-P_{\mathcal{F}} Z$, and replace $X$ by $P_{\mathcal{D}} Y$ as given by its expression on the right-hand-side of (15). The condition to determine $\lambda$ then becomes

$$
\lambda\left\langle\left\langle P_{\mathcal{D}}\left(P_{\mathcal{F}} Z\right)-P_{\mathcal{D}} Z, P_{\mathcal{F}} Z-Z\right\rangle\right\rangle+\left\langle\left\langle P_{\mathcal{D}} Z-P_{\mathcal{F}} Z, P_{\mathcal{F}} Z-Z\right\rangle\right\rangle=0 .
$$

The first inner-product term on the left-hand-side may be written as

$$
\left\langle\left\langle P_{\mathcal{D}}\left(P_{\mathcal{F}} Z\right)-P_{\mathcal{D}} Z, P_{\mathcal{D}}\left(P_{\mathcal{F}} Z\right)-Z\right\rangle\right\rangle,
$$

due to the fact that $P_{\mathcal{D}}$ is an orthogonal projection onto $\mathcal{D}$. Hence,

$$
\lambda=\frac{\left\langle\left\langle P_{\mathcal{D}} Z-P_{\mathcal{F}} Z, Z-P_{\mathcal{F}} Z\right\rangle\right\rangle}{\left\langle\left\langle P_{\mathcal{D}}\left(P_{\mathcal{F}} Z\right)-P_{\mathcal{D}} Z, P_{\mathcal{D}}\left(P_{\mathcal{F}} Z\right)-Z\right\rangle\right\rangle} .
$$

The expression (14) is then obtained by writing in the denominator

$$
P_{\mathcal{D}}\left(P_{\mathcal{F}} Z\right)-Z=P_{\mathcal{D}}\left(P_{\mathcal{F}} Z\right)-P_{\mathcal{D}} Z+P_{\mathcal{D}} Z-Z
$$

and computing the inner product, using the fact that $P_{\mathcal{D}} Z-Z$ is orthogonal onto $\mathcal{D}$.

Conversely, when for a given $Z \in \mathcal{H}, Z \neq P_{\mathcal{F}} Z$, the number $\lambda$ given by (14) exists, it follows by a routine control that, when that value of $\lambda$ is substituted into (13), then $P_{\mathcal{D}} Y \in \mathcal{P}$.

In the special case that the point $Z$ belongs to $\mathcal{D}$ we obtain the following interesting result about the corresponding value of $\lambda$. Denoting now $Z$ as $X$, we have $P_{\mathcal{D}} X=X, P_{\mathcal{F}} X \neq X$, and

$$
\begin{equation*}
\lambda=\frac{\left\|X-P_{\mathcal{F}} X\right\|^{2}}{\left\|P_{\mathcal{D}}\left(P_{\mathcal{F}} X\right)-X\right\|^{2}}=1+\frac{\left\|P_{\mathcal{F}} X-P_{\mathcal{D}}\left(P_{\mathcal{F}} X\right)\right\|^{2}}{\left\|P_{D}\left(P_{\mathcal{F}} X\right)-X\right\|^{2}} \tag{16}
\end{equation*}
$$

where the last equality is obtained by using the Pythagorean theorem on the triangle with vertices $X, P_{\mathcal{F}} X, P_{\mathcal{D}}\left(P_{\mathcal{F}} X\right)$. Hence, when $X \in D$, $X \notin \mathcal{F}$, and the corresponding value of $\lambda$ exists, then this value is never smaller than 1 .

Lemma 3. When $X \in \mathcal{D}, X \notin \mathcal{F}$, then $\lambda$ is well-defined.
Proof. We have to show that the denominator appearing in the expressions of (16) can not be zero. To this end, use is made of the following formulas concerning distances in H and $\mathcal{H}$ related to a point $X \equiv q(\boldsymbol{x})$ in $\mathcal{D} \subset \mathcal{H}$ with $\boldsymbol{x} \in H:$

$$
\begin{align*}
\left\|X-P_{\mathcal{F}} X\right\|^{2} & =\sum_{j=1}^{r}\left|\boldsymbol{x}-P_{j} \boldsymbol{x}\right|^{2},  \tag{17}\\
\left\|X-P_{\mathcal{D}}\left(P_{\mathcal{F}} X\right)\right\|^{2} & =r\left|\sum_{j=1}^{r} \frac{1}{r}\left(P_{j} \boldsymbol{x}-\boldsymbol{x}\right)\right|^{2} . \tag{18}
\end{align*}
$$

For $\boldsymbol{c} \in C^{*} \equiv \bigcap_{j=1}^{r} C_{j} \subset H$ we have
$\left\langle\boldsymbol{x}-\boldsymbol{c}, \sum_{j=1}^{r} \frac{1}{r}\left(P_{j} \boldsymbol{x}-\boldsymbol{x}\right)\right\rangle=\sum_{j=1}^{r} \frac{1}{r}\left\langle\boldsymbol{x}-P_{j} \boldsymbol{x}, P_{j} \boldsymbol{x}-\boldsymbol{x}\right\rangle+\sum_{j=1}^{r} \frac{1}{r}\left\langle P_{j} \boldsymbol{x}-\boldsymbol{c}, P_{j} \boldsymbol{x}-\boldsymbol{x}\right\rangle$,
and each term in the last sum is non-positive due to a well-known property of projections. Hence,

$$
\begin{equation*}
\left\langle\boldsymbol{x}-\boldsymbol{c}, \sum_{j=1}^{r} \frac{1}{r}\left(P_{j} \boldsymbol{x}-\boldsymbol{x}\right)\right\rangle \leq-\frac{1}{r} \sum_{j=1}^{r}\left|\boldsymbol{x}-P_{j} \boldsymbol{x}\right|^{2} . \tag{19}
\end{equation*}
$$

Assuming now that the denominator of (16) is zero, it follows from (18) that $\sum_{j=1}^{r} \frac{1}{r}\left(P_{j} \boldsymbol{x}-\boldsymbol{x}\right)=0$, and from (19) it follows also that $\sum_{j=1}^{r} \mid \boldsymbol{x}-$ $\left.P_{j} \boldsymbol{x}\right|^{2}=0$; using (17) we conclude that then $X=P_{\mathcal{F}} X$, contrary to our assumptions.

The foregoing result is interesting for our algorithm in the sense that, when an iteration point $X_{k}$ in $\mathcal{D}$ has been obtained that not yet belongs to $\mathcal{F}$, and when that iteration point $X_{k}$ is used as the starting point for the next iteration step in $\mathcal{H}$, then we are sure that the corresponding $\lambda$ is well-defined; moreover, it is not smaller than 1.

Lemma 4. Let $X \in \mathcal{D}, X \notin \mathcal{F}$, and put $Y=X+\lambda\left(P_{\mathcal{F}} X-X\right)$ with $\lambda$ as given by (16). Then $\lambda=1$ if and only if $Y \in \mathcal{D}$.

Proof. Suppose $\lambda=1$. Then $Y=P_{\mathcal{F}} X$ and $P_{\mathcal{D}} Y=P_{\mathcal{D}}\left(P_{\mathcal{F}} X\right)$. Moreover, from the last expression in (16) we also conclude that then $P_{\mathcal{D}}\left(P_{\mathcal{F}} X\right)=P_{\mathcal{F}} X$. Hence, $Y=P_{\mathcal{D}} Y$, and so $Y \in \mathcal{D}$.

Conversely, suppose that $Y \in \mathcal{D}$, which means that $P_{\mathcal{D}} Y=Y$. Equating the expressions of $Y$ and $P_{\mathcal{D}} Y$ leads to $\lambda\left(P_{\mathcal{F}} X-P_{\mathcal{D}}\left(P_{\mathcal{F}} X\right)\right)=0$. As $\lambda \geq 1$, there results that $P_{\mathcal{F}} X=P_{\mathcal{D}}\left(P_{\mathcal{F}} X\right)$ and, again from (16), we conclude that $\lambda=1$.

In particular, when $X \in \mathcal{D}, X \notin \mathcal{F}$, and $\lambda=1$, then $P_{\mathcal{D}} Y=P_{\mathcal{F}} X \in$ $\mathcal{F} \cap \mathcal{D}$. For our algorithm this means that, when taking an iteration point $X_{k}$ in $\mathcal{D}$ as the starting point for the next iteration step in $\mathcal{H}$, and when $\lambda$ corresponding to $X_{k}$ has the value 1 , then $X_{k+1}$ will belong to $\mathcal{F} \cap \mathcal{D}$ when $P_{\mathcal{D}} Y$ is taken as the next iteration point $X_{k+1}$.

For a point $Z \in \mathcal{H}$ with $Z \notin \mathcal{D}$ and $Z \notin \mathcal{F}$, Lemma 4 can be adapted as follows.

Lemma 5. Let $Z \notin \mathcal{D}, Z \notin \mathcal{F}$, and put $Y=Z+\lambda\left(P_{\mathcal{F}} Z-Z\right)$ with $\lambda$ as given by (14). When $Y \in \mathcal{D}$, then necessarily $\lambda=1$.

Proof. By assumption, $P_{\mathcal{D}} Y=Y$ and $P_{\mathcal{D}} Y \in \mathcal{P}$. Hence,

$$
\begin{aligned}
0 & =\left\langle\left\langle Z-P_{\mathcal{F}} Z, Y-P_{\mathcal{F}} Z\right\rangle\right\rangle=\left\|Z-P_{\mathcal{F}} Z\right\|^{2}+\lambda\left\langle\left\langle Z-P_{\mathcal{F}} Z, P_{\mathcal{F}} Z-Z\right\rangle\right\rangle \\
& =(1-\lambda)\left\|Z-P_{\mathcal{F}} Z\right\|^{2},
\end{aligned}
$$

which leads to the result, as $Z \neq P_{\mathcal{F}} Z$.
Let $Z_{0}$ be an element of $\mathcal{H}$, whether or not in $\mathcal{D}$, but $Z_{0} \notin \mathcal{F}$. Put $Y_{1}=Z_{0}+\lambda_{1}\left(P_{\mathcal{F}} Z_{0}-Z_{0}\right)$, and suppose that the corresponding $\lambda_{1}$-value in order that $P_{\mathcal{D}} Y_{1}$ should belong to the supporting hyperplane $\mathcal{P}_{1}$ at $P_{\mathcal{F}} Z_{0}$, exists and is bigger than 1. Denote $P_{\mathcal{D}} Y_{1}$ as $X_{1}$, and put, for some $\gamma>0$

$$
\begin{equation*}
Z_{1}=Y_{1}+(1+\gamma)\left(X_{1}-Y_{1}\right) \tag{20}
\end{equation*}
$$

In view of Lemmas 4 and 5 we know that $X_{1} \neq Y_{1}$. Hence, $Z_{1}$ is the result of a relaxed projection of $Y_{1}$ onto D . We also note that $P_{\mathcal{D}} Z_{1}=X_{1}$.

We are interested in the fact whether or not $Z_{1}$ can belong to $\mathcal{F}$.
Lemma 6. When $Z_{0} \notin F, \lambda_{1}>1, \gamma>0$, then $Z_{1} \notin \mathcal{F}$.
Proof. Replacing in the following inner product $Z_{1}$ by its expression given in (20) we obtain

$$
\begin{aligned}
\left\langle\left\langle Z_{0}-P_{\mathcal{F}} Z_{0}, Z_{1}-P_{\mathcal{F}} Z_{0}\right\rangle\right\rangle= & \left\langle\left\langle Z_{0}-P_{\mathcal{F}} Z_{0}, X_{1}-P_{\mathcal{F}} Z_{0}\right\rangle\right\rangle \\
& +\gamma\left\langle\left\langle Z_{0}-P_{\mathcal{F}} Z_{0}, X_{1}-Y_{1}\right\rangle\right\rangle .
\end{aligned}
$$

The first inner product on the right-hand-side is zero due to the definition of $\lambda_{1}$. Replacing in the second inner product $Y_{1}$ by its expression in $Z_{0}$ and $P_{\mathcal{F}} Z_{0}$ we obtain

$$
\left\langle\left\langle Z_{0}-P_{\mathcal{F}} Z_{0}, Z_{1}-P_{\mathcal{F}} Z_{0}\right\rangle\right\rangle=\gamma \lambda_{1}\left\|Z_{0}-P_{\mathcal{F}} Z_{0}\right\|^{2}+\gamma\left\langle\left\langle Z_{0}-P_{\mathcal{F}} Z_{0}, X_{1}-Z_{0}\right\rangle\right\rangle .
$$

To obtain the inner product on the right-hand-side we write

$$
\begin{gathered}
\left\|X_{1}-P_{\mathcal{F}} Z_{0}\right\|^{2}=\left\|X_{1}-Z_{0}\right\|^{2}+\left\|Z_{0}-P_{\mathcal{F}} Z_{0}\right\|^{2}+2\left\langle\left\langle X_{1}-Z_{0}, Z_{0}-P_{\mathcal{F}} Z_{0}\right\rangle\right\rangle \\
\quad=\left\|X_{1}-P_{\mathcal{F}} Z_{0}\right\|^{2}+2\left\|Z_{0}-P_{\mathcal{F}} Z_{0}\right\|^{2}+2\left\langle\left\langle X_{1}-Z_{0}, Z_{0}-P_{\mathcal{F}} Z_{0}\right\rangle\right\rangle
\end{gathered}
$$

due to the Pythagorean theorem for the right triangle with vertices $Z_{0}$, $P_{\mathcal{F}} Z_{0}$ and $X_{1}$. Hence,

$$
\left\langle\left\langle X_{1}-Z_{0}, Z_{0}-P_{\mathcal{F}} Z_{0}\right\rangle\right\rangle=-\left\|Z_{0}-P_{\mathcal{F}} Z_{0}\right\|^{2} .
$$

Using this result in the former equality we obtain

$$
\left\langle\left\langle Z_{0}-P_{\mathcal{F}} Z_{0}, Z_{1}-P_{\mathcal{F}} Z_{0}\right\rangle\right\rangle=\gamma\left(\lambda_{1}-1\right)\left\|Z_{0}-P_{\mathcal{F}} Z_{0}\right\|^{2} .
$$

Hence, $\left\langle\left\langle Z_{0}-P_{\mathcal{F}} Z_{0}, Z_{1}-P_{\mathcal{F}} Z_{0}\right\rangle\right\rangle$ is strictly positive, from which we conclude, due to a property of projections, that $Z_{1} \notin \mathcal{F}$.

Now we are ready to describe our algorithm to construct the approximating sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$ in $\mathcal{D}$ in order to obtain a point in $\mathcal{D} \cap \mathcal{F}$.

## Algorithm

Starting from some point $Z_{0} \equiv X_{0} \in \mathcal{D}$, suppose that points $X_{1}, \ldots$ $\ldots, X_{k}$ in $\mathcal{D}$ and $Z_{1}, \ldots, Z_{k}$ in $\mathcal{H}$ (whether or not in $\mathcal{D}$ ) have been obtained, with $P_{\mathcal{D}} Z_{k}=X_{k}, X_{k} \notin \mathcal{F} \cap \mathcal{D}, Z_{k} \notin \mathcal{F}$. Then:
(i) IF $P_{\mathcal{D}}\left(P_{\mathcal{F}} Z_{k}\right)=X_{k}$, then take $X_{k}$ as new $Z_{k}$ and go to (ii)
(ii) ELSE, compute $\lambda_{k+1}$ by (14)(with $Z$ replaced by $Z_{k}$ )
a. IF $\lambda_{k+1}>1$, let

$$
\begin{aligned}
Y_{k+1} & =Z_{k}+\lambda_{k+1}\left(P_{\mathcal{F}} Z_{k}-Z_{k}\right), \\
X_{k+1} & =P_{\mathcal{D}} Y_{k+1}, \\
Z_{k+1} & =Y_{k+1}+\left(1+\gamma_{k+1}\right)\left(X_{k+1}-Y_{k+1}\right),
\end{aligned}
$$

(where $\gamma_{k+1}>0$ will be determined further on)
b. ELSE, let $X_{k+1}=P_{\mathcal{D}}\left(P_{\mathcal{F}} Z_{k}\right)$,

$$
Z_{k+1}=X_{k+1} .
$$

As has been shown in Lemma 3 and in the remark following (16), whenever $Z_{k} \equiv X_{k} \in \mathcal{D}$ but $X_{k} \notin \mathcal{F}$ then $\lambda_{k+1}$ is well-defined and is never smaller than 1. Hence, in that case the steps (i) and (ii)a. of the algorithm need not be investigated; moreover, when $\lambda_{k+1}=1$, then $X_{k+1} \in \mathcal{D} \cap \mathcal{F}$ according to Lemma 4.

Concerning a suitable value for $\gamma_{k+1}$ needed in (ii)b., we remember that taking $Z_{k+1}$ instead of $X_{k+1}$ as a new starting point in $\mathcal{H}$ for the computation of the next update $X_{k+2}$ has the purpose of possibly interrupting, at some steps, the monotone behaviour of the approximating sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$; this needs a $\gamma_{k+1}$ that is different from zero; we choose it strictly positive. On the other hand, however, we need convergence of the sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$, and so we want an interruption that somehow takes into consideration part of the approximation already obtained. These facts seem to imply that, in any case, the distance between $X_{k+1}$ and $Z_{k+1}$ should tend to zero with growing $k$. To this end, for $M$ a fixed positive number we let $\gamma_{k+1}$ be dependent on $\frac{M}{k+1}$. In order to take care of the fact that $\left\{\lambda_{k+1}\right\}_{k=0}^{+\infty}$ may be unbounded, and hence also $\left\{\left\|X_{k+1}-Y_{k+1}\right\|\right\}_{k=0}^{+\infty}$, let $B$ be a (big) fixed positive number; we then put

$$
\begin{gathered}
\alpha_{k+1}=\min \left(\frac{1}{\lambda_{k+1}}, \frac{M}{k+1}\right), \quad \beta_{k+1}=\min \left(1, \frac{B}{\left\|X_{k+1} Y_{k+1}\right\|}\right), \\
\gamma_{k+1}=\alpha_{k+1} \beta_{k+1} .
\end{gathered}
$$

## 3. Convergence of the constructed sequence

In this section we show that the sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$ in $\mathcal{D}$ is convergent to a point of $\mathcal{D} \cap \mathcal{F}$. Due to the fact that the construction of the points $X_{k+1}$ in $\mathcal{D}$ is different according to the corresponding values of $\lambda_{k+1}$, we have to consider separately the properties of elements of the sequence corresponding to values of $\lambda_{k+1}$ that are bigger than 1 , or not.

## A. Properties corresponding to $\lambda_{k+1}>1$

Let us suppose that we obtained $X_{k}$ in $\mathcal{D}, Z_{k}$ in $\mathcal{H}$ with $P_{\mathcal{D}}\left(Z_{k}\right)=X_{k}$ (it is possible that $X_{k}=Z_{k}$ ), and that $\lambda_{k+1}$, given by (14) with $Z$ replaced by $Z_{k}$, is well-defined and is bigger than 1 . According to (ii)b. of the algorithm, we have

$$
\begin{equation*}
X_{k+1}=P_{D} Y_{k+1}=X_{k}+\lambda_{k+1}\left(P_{\mathcal{D}}\left(P_{\mathcal{F}} Z_{k}\right)-X_{k}\right) \tag{21}
\end{equation*}
$$

and $X_{k+1}$ belongs to the hyperplane $\mathcal{P}_{k+1}$ supporting $\mathcal{F}$ at $P_{\mathcal{F}} Z_{k}$; its equation is

$$
\begin{equation*}
\mathcal{P}_{k+1}=\left\{X \in \mathcal{H}:\left\langle\left\langle X-P_{\mathcal{F}} Z_{k}, P_{\mathcal{F}} Z_{k}-Z_{k}\right\rangle\right\rangle=0\right\} . \tag{22}
\end{equation*}
$$

## Lemma 7.

$$
\begin{equation*}
\left\langle\left\langle X_{k}-X_{k+1}, V-X_{k+1}\right\rangle\right\rangle \leq 0, \quad \forall V \in \mathcal{D} \cap \mathcal{F} . \tag{23}
\end{equation*}
$$

Proof. Using the facts that $X_{k}-Z_{k}$ and $Y_{k+1}-X_{k+1}$ are orthogonal onto $D$, and that $V-X_{k+1} \in \mathcal{D}$, we have

$$
\begin{aligned}
\left\langle\left\langle X_{k}-X_{k+1}, V-X_{k+1}\right\rangle\right\rangle & =\left\langle\left\langle Z_{k}-X_{k+1}, V-X_{k+1}\right\rangle\right\rangle \\
& =\left\langle\left\langle Z_{k}-Y_{k+1}, V-X_{k+1}\right\rangle\right\rangle .
\end{aligned}
$$

Replacing $Y_{k+1}$ by its expression in $Z_{k}$ and $P_{\mathcal{F}} Z_{k}$ leads to

$$
\begin{aligned}
& \left\langle\left\langle X_{k}-X_{k+1}, V-X_{k+1}\right\rangle\right\rangle=-\lambda_{k+1}\left\langle\left\langle P_{\mathcal{F}} Z_{k}-Z_{k}, V-X_{k+1}\right\rangle\right\rangle \\
& \left.\quad=\lambda_{k+1}\left\langle\left\langle Z_{k}-P_{\mathcal{F}} Z_{k}, V-P_{\mathcal{F}} Z_{k}\right\rangle\right\rangle+\lambda_{k+1}\left\langle Z_{k}-P_{\mathcal{F}} Z_{k}, P_{\mathcal{F}} Z_{k}-X_{k+1}\right\rangle\right\rangle .
\end{aligned}
$$

In the right-hand-side, the last term is zero because $X_{k+1} \in \mathcal{P}_{k+1}$, while the first term is nonpositive due to a well-known property of projections.

Hence the result follows.

For points $V \in \mathcal{D} \cap \mathcal{F}$ we further have

$$
\left\|X_{k}-V\right\|^{2}=\left\|X_{k}-X_{k+1}\right\|^{2}+\left\|X_{k+1}-V\right\|^{2}+2\left\langle\left\langle X_{k}-X_{k+1}, X_{k+1}-V\right\rangle\right\rangle
$$

and taking (23) into account we derive that

$$
\begin{equation*}
\left\|X_{k}-V\right\|^{2} \geq\left\|X_{k}-X_{k+1}\right\|^{2}+\left\|X_{k+1}-V\right\|^{2}, \quad \forall V \in \mathcal{D} \cap \mathcal{F} \tag{24}
\end{equation*}
$$

In particular, we conclude that, when $X_{k+1}$ is obtained using a value of $\lambda_{k+1}$ that is bigger than 1, then

$$
\left\|X_{k}-V\right\| \geq\left\|X_{k+1}-V\right\|, \quad \forall V \in \mathcal{D} \cap \mathcal{F}
$$

which expresses the monotone approach to $D \cap F$ in that case.
Let us suppose now that, except at a finite number of steps, all values of $\lambda_{k+1}$ are bigger than 1 ; let us restrict to such values of $k$. Then from (24) we derive that

$$
\begin{equation*}
\left\|X_{k}-X_{k+1}\right\| \rightarrow 0 \quad \text { when } \quad k \rightarrow+\infty \tag{25}
\end{equation*}
$$

Making use of the Pythagorean theorem for the right triangles with vertices $Z_{k}, X_{k}, X_{k+1}$ and $Z_{k}, P_{\mathcal{F}} Z_{k}, X_{k+1}$ respectively, we obtain

$$
\begin{aligned}
\left\|X_{k+1}-X_{k}\right\|^{2} & =\left\|Z_{k}-X_{k+1}\right\|^{2}-\left\|Z_{k}-X_{k}\right\|^{2} \\
& =\left\|Z_{k}-P_{\mathcal{F}} Z_{k}\right\|^{2}+\left\|P_{\mathcal{F}} Z_{k}-X_{k+1}\right\|^{2}-\left\|Z_{k}-X_{k}\right\|^{2} .
\end{aligned}
$$

Due to our choice of $\gamma_{k+1}$ in the algorithm we know that $\left\|Z_{k}-X_{k}\right\| \rightarrow 0$ when $k \rightarrow+\infty$. Hence, we obtain that also $\left\|Z_{k}-P_{\mathcal{F}} Z_{k}\right\| \rightarrow 0$ and $\left\|P_{\mathcal{F}} Z_{k}-X_{k+1}\right\| \rightarrow 0$ when $k \rightarrow+\infty$. As

$$
\left\|X_{k}-P_{\mathcal{F}} X_{k}\right\| \leq\left\|X_{k}-P_{\mathcal{F}} Z_{k}\right\| \leq\left\|X_{k}-X_{k+1}\right\|+\left\|X_{k+1}-P_{\mathcal{F}} Z_{k}\right\|
$$

there results that also

$$
\begin{equation*}
\left\|X_{k}-P_{\mathcal{F}} X_{k}\right\| \rightarrow 0 \quad \text { when } \quad k \rightarrow+\infty \tag{26}
\end{equation*}
$$

Combining the fact that the sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$ is bounded with the relations (24), (25) and (26), we have the necessary information needed to prove that the sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$ is weakly (and hence in norm) convergent to a point of $\mathcal{F} \cap \mathcal{D}$. We do not represent the proof here, but refer instead to [6] where it has been shown (in another situation) that from the mentioned properties convergence may be obtained.

## B. Properties corresponding to $\lambda_{k+1} \leq 1$

We know that, when $Z_{k} \equiv X_{k} \in \mathcal{D}$, then $\lambda_{k+1} \geq 1$, and when $\lambda_{k+1}=1$ then $X_{k+1} \equiv P_{\mathcal{F}} Z_{k}$ is a point of $\mathcal{D} \cap \mathcal{F}$ and the procedure is finished; so, when $Z_{k} \equiv X_{k}$ and the next $X_{k+1}$ is not yet in $\mathcal{D} \cap \mathcal{F}$, we always have that $\lambda_{k+1}>1$; we also started with $Z_{0} \equiv X_{0} \in \mathcal{D}$.

On the other hand, when $Z_{k} \neq X_{k}$ and $\lambda_{k+1} \leq 1$, then we know from (ii)c. of the algorithm that $X_{k+1}=P_{\mathcal{D}}\left(P_{\mathcal{F}} Z_{k}\right)=Z_{k+1}$. It is at such point $X_{k+1}$ that the monotone behaviour of the sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$ may be (but is not necessarily) interrupted, i.e., it is not necessarily true that $\left\|X_{k+1}-V\right\| \leq\left\|X_{k}-V\right\|, \forall V \in \mathcal{F} \cap \mathcal{D}$, but, as $Z_{k+1} \equiv X_{k+1} \in \mathcal{D}$, we always have that $\lambda_{k+2} \geq 1$.

Hence, in the case that we are investigating now, we may suppose that the sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$ has an (infinite) subsequence $\left\{X_{n_{j}}\right\}_{j=0}^{+\infty}$ with $X_{n_{0}} \equiv Z_{n_{0}}, X_{n_{1}} \equiv Z_{n_{1}}, \ldots, X_{n_{j}} \equiv Z_{n_{j}}, \ldots$, with $n_{0}=0$, and those points were obtained from a corresponding $\lambda_{n_{j}}$-value not bigger than 1 (except for the starting point $X_{0}$ ). The $\lambda$-values for points between two consecutive points $X_{n_{p}}$ and $X_{n_{q}}$ of the subsequence are strictly bigger than 1; i.e., we have that $\lambda_{n_{p}} \leq 1$ (except for $n_{p}=0$ ), $\lambda_{n_{p}+1}>1, \ldots, \lambda_{n_{q}-1}>1, \lambda_{n_{q}} \leq 1$.

We show that, due to our choice of $\gamma_{k+1}$, a weaker form of monotony is still available, which will be sufficient to guarantee convergence of the sequence. We first prove the following lemma, valid for $\lambda$-values that are strictly bigger than 1 .

Lemma 8. Suppose, with the usual notations, that for some $k \in \mathbb{Z}^{+}$ we have that $X_{k} \in \mathcal{D}, Z_{k} \notin \mathcal{D}, \lambda_{k+1}>1$, giving rise to $Y_{k+1}, X_{k+1}$ and $Z_{k+1}$. Then

$$
\begin{equation*}
\left\|Z_{k+1}-X_{k+1}\right\|^{2} \leq\left\|Z_{k}-X_{k}\right\|^{2}+\left\|X_{k}-X_{k+1}\right\|^{2} \tag{27}
\end{equation*}
$$

Proof. The straight line through $P_{\mathcal{F}} Z_{k}$ that is parallel to the straight line through $Y_{k+1}$ and $Z_{k+1}$, intersects the straight line $Z_{k} X_{k+1}$ in a point $S_{k}$. From the similarity of the triangles with vertices $Z_{k}, P_{\mathcal{F}} Z_{k}, S_{k}$ and $Z_{k}, Y_{k+1}, X_{k+1}$ we deduce that $\left\|Y_{k+1}-X_{k+1}\right\|=\lambda_{k+1}\left\|P_{\mathcal{F}} Z_{k}-S_{k}\right\|$. But $Z_{k+1}=X_{k+1}+\gamma_{k+1}\left(X_{k+1}-Y_{k+1}\right)$, and so

$$
\left\|Z_{k+1}-X_{k+1}\right\| \leq \gamma_{k+1} \lambda_{k+1}\left\|P_{\mathcal{F}} Z_{k}-S_{k}\right\| \leq\left\|P_{\mathcal{F}} Z_{k}-S_{k}\right\|
$$

Moreover, $\left\|P_{\mathcal{F}} Z_{k}-S_{k}\right\| \leq\left\|Z_{k}-X_{k+1}\right\|$, which follows from inspection of the triangle with vertices $Z_{k}, P_{\mathcal{F}} Z_{k}, X_{k+1}$. Finally, $Z_{k}-X_{k}$ is orthogonal onto $X_{k}-X_{k+1}$, and so we obtain

$$
\left\|Z_{k+1}-X_{k+1}\right\|^{2} \leq\left\|Z_{k}-X_{k+1}\right\|^{2}=\left\|Z_{k}-X_{k}\right\|^{2}+\left\|X_{k}-X_{k+1}\right\|^{2} .
$$

Using (27) and the Fejér-monotony property (24) a number of times leads to the following weaker form of monotony in the case where an interrupting subsequence $\left\{X_{n_{j}}\right\}_{j=0}^{+\infty}$ may exist.

Lemma 9. Suppose that $X_{n_{p}} \equiv Z_{n_{p}}$ and $X_{n_{q}} \equiv Z_{n_{q}}$ are consecutive points of the subsequence $\left\{X_{n_{j}}\right\}_{j=0}^{+\infty}$. Then, $\forall V \in \mathcal{D} \cap \mathcal{F}$ we have

$$
\begin{equation*}
\left\|X_{n_{p}}-V\right\|^{2} \geq\left\|X_{n_{q}}-V\right\|^{2}+\frac{3}{4}\left\|X_{n_{p}}-X_{n_{p}+1}\right\|^{2} \tag{28}
\end{equation*}
$$

Proof. In order not to overload the proof with cumbersome indexes, we present a proof with concrete small numbers. So, let $n_{p}=0, n_{q}=4$, $V \in \mathcal{F} \cap \mathcal{D}$. Then we first have

$$
\begin{equation*}
\left\|V-X_{4}\right\|^{2} \leq\left\|V-P_{\mathcal{F}} Z_{3}\right\|^{2} \leq\left\|V-Z_{3}\right\|^{2}=\left\|V-X_{3}\right\|^{2}+\left\|X_{3}-Z_{3}\right\|^{2} \tag{29}
\end{equation*}
$$

by using the non-expansivity property of projections and the fact that $Z_{3}-X_{3}$ is orthogonal onto $\mathcal{D}$. Using (27) for $k=2$ leads to

$$
\left\|V-X_{4}\right\|^{2} \leq\left\|V-X_{3}\right\|^{2}+\left\|Z_{2}-X_{2}\right\|^{2}+\left\|X_{2}-X_{3}\right\|^{2} .
$$

As $\lambda_{3}>1$, inequality (24) may be used for $k=2$, giving

$$
\begin{gathered}
\left\|V-X_{4}\right\|^{2} \leq\left\|V-X_{2}\right\|^{2}-\left\|X_{2}-X_{3}\right\|^{2}+\left\|Z_{2}-X_{2}\right\|^{2}+\left\|X_{2}-X_{3}\right\|^{2} \\
=\left\|V-X_{2}\right\|^{2}+\left\|X_{2}-Z_{2}\right\|^{2},
\end{gathered}
$$

and it is instructive to compare this result with (29). Hence, again applying (24) and (27), now for $k=1$, leads to

$$
\begin{gathered}
\left\|V-X_{4}\right\|^{2} \leq\left\|V-X_{1}\right\|^{2}-\left\|X_{1}-X_{2}\right\|^{2}+\left\|Z_{1}-X_{1}\right\|^{2}+\left\|X_{1}-X_{2}\right\|^{2} \\
=\left\|V-X_{1}\right\|^{2}+\left\|Z_{1}-X_{1}\right\|^{2} .
\end{gathered}
$$

For the final step we use two different things. First, according to (24) we have that $\left\|V-X_{1}\right\|^{2} \leq\left\|V-X_{0}\right\|^{2}-\left\|X_{0}-X_{1}\right\|^{2}$. For the term $\left\|Z_{1}-X_{1}\right\|^{2}$,
however, we remark that now $X_{0} \in \mathcal{D}$ (remember that this is true for the general index $n_{p}$ ). Hence, the triangle with vertices $X_{0}, P_{\mathcal{F}} X_{0}$ and $X_{1}$ has a right angle in $P_{\mathcal{F}} X_{0}$. From some elementary geometry it follows then easily that

$$
\left\|Z_{1}-X_{1}\right\|^{2} \leq\left\|P_{\mathcal{F}} X_{0}-P_{\mathcal{D}}\left(P_{\mathcal{F}} X_{0}\right)\right\|^{2} \leq \frac{1}{4}\left\|X_{0}-X_{1}\right\|^{2}
$$

So, we finally obtain

$$
\left\|V-X_{4}\right\|^{2} \leq\left\|V-X_{0}\right\|^{2}-\left\|X_{0}-X_{1}\right\|^{2}+\frac{1}{4}\left\|X_{0}-X_{1}\right\|^{2},
$$

which is precisely formula (28) in our case.
Inequality (28) gives an idea about the measure of interrupting the monotone behaviour at the points of the subsequence: although it may be true that $\left\|X_{n_{q}}-W\right\|>\left\|X_{n_{q}-1}-W\right\|$ for some points $W$ in $\mathcal{F} \cap \mathcal{D}$, the monotony is repaired with respect to the foregoing interruption point.

Due to inequalities (28) and (24) we are able to prove that, also in case that the sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$ has an infinite interrupting subsequence $\left\{X_{n_{p}}\right\}_{p=0}^{+\infty}$, the whole sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$ is convergent to a point $A \in \mathcal{D} \cap \mathcal{F}$. Again, the method described in [6] may be used. So we just sketch the main ideas. The subsequence $\left\{X_{n_{p}}\right\}_{p=0}^{+\infty}$ has a subsequence, say $\left\{X_{q}^{\prime}\right\}_{q=0}^{+\infty}$, that weakly converges to a point $A \in \mathcal{D}$; this point also belongs to $\mathcal{F}$, as $\left\|P_{\mathcal{F}} X_{n_{p}}-X_{n_{p}}\right\| \rightarrow 0$. Moreover, each subsequence of $\left\{X_{n_{p}}\right\}_{p=0}^{+\infty}$ that weakly converges has the same point $A$ as its weak limit; hence, $\left\{X_{n_{p}}\right\}_{p=0}^{+\infty}$ is (weakly and in) norm convergent to $A \in \mathcal{D} \cap \mathcal{F}$. Convergence of the whole sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$ to $A$ then follows from the fact that, for $k>n_{q}$ we have that $\left\|X_{k}-A\right\| \leq\left\|X_{n_{q}}-A\right\|$.

We summarise the results in the following theorem.
Theorem. Suppose that in $\mathbb{R}^{m}$, $r$ closed convex sets $\left\{C_{j}\right\}_{j=1}^{r}$ with nonempty intersection $\bigcap_{j=1}^{r} C_{j}$ are given. Let $q$ be the natural imbedding of $\mathbb{R}^{m}$ in $\left(\mathbb{R}^{m}\right)^{r}$, and put $\mathcal{D}=q\left(\mathbb{R}^{m}\right), \mathcal{F}=C_{1} \times C_{2} \times \cdots \times C_{r}$. For a given point $\boldsymbol{x}_{0} \in \mathbb{R}^{m}$, put $X_{0} \equiv q\left(\boldsymbol{x}_{0}\right)$. Then the sequence $\left\{X_{k}\right\}_{k=0}^{+\infty}$, constructed according to the algorithm mentioned in Section 2, is convergent to a point of $\mathcal{F} \cap \mathcal{D}$. Hence, the sequence $\left\{\boldsymbol{x}_{k}\right\}_{k=0}^{+\infty}$ in $\mathbb{R}^{m}$ with $\boldsymbol{x}_{k}=q^{-1}\left(X_{k}\right)$ is convergent to a point in $\bigcap_{j=1}^{r} C_{j}$.

## 4. Examples and concluding remarks

In this last section we illustrate the algorithm given in Section 2, and we compare the results with the ones corresponding to two other methods. The following algorithms have been used:
PP: The method of pure projections in a sequential manner: when $\boldsymbol{x}_{k}$ is the current iteration point, $\boldsymbol{x}_{k+1}$ is obtained by $\boldsymbol{x}_{k+1}=T \boldsymbol{x}_{k}$, with $T \equiv P_{r} \ldots P_{2} P_{1}$, and with $P_{j}$ denoting the projection operator onto $C_{j}$.
PAR: The parallel projection given in (2) with fixed equal weights (i.e., $\mu_{k+1}(j)=\frac{1}{r}$ for each k ), and with $\lambda_{k+1}$ determined as in [6], i.e.,

$$
\lambda_{k+1}=\frac{\sum_{j=1}^{r} \frac{1}{r}\left|\boldsymbol{x}_{k}-P_{j} \boldsymbol{x}_{k}\right|^{2}}{\left|\boldsymbol{x}_{k}-\sum_{j=1}^{r} \frac{1}{r} P_{j} \boldsymbol{x}_{k}\right|^{2}}
$$

NEW0 and NEW3: The algorithm presented in this paper, with $M=1$ and $M=10^{3}$ respectively, and with $B=10^{6}$ in both cases.

As a first example, we take as closed convex sets the twelve disks $\left\{C_{j}\right\}_{j=1}^{12}$ mentioned in formula (1), Section 1 ; their intersection (in fact determined by $C_{1}$ and $C_{12}$ ) contains more than one point. Starting from some given point in the plane, we want to obtain a point in their intersection; explicit expressions for $P_{j}$ may be found in [2].

In Table 1 we have mentioned, for the given eight starting points $((-3,0), \ldots)$, either the number of iterations needed to obtain a point in the intersection (this is a positive integer), or the sum of the distances of the current iteration point to the twelve sets $C_{j}$, after 25 and 50 iterations respectively. The results in Table 1 suggest that the convergence behaviour of the algorithm presented in this paper is much less dependent on the starting point than in the other cases; extremely slow convergence did not appear, for a lot of starting points.

In our second example we considered in $\mathbb{R}^{3}$ with variables $x, y, z$ the intersection of the ball given by $x^{2}+y^{2}+z^{2} \leq R^{2}$ (with $\frac{1}{6} \leq R \leq 1$ ), and the three half-spaces given by $x+y+4 z \leq 1, x+y-4 z \leq 1,-x+y-8 z \leq 1$. The four algorithms mentioned in example 1 all gave rather comparable results: a point in the intersection was obtained in at most 4 iterations, for different starting points and for different values of the radius $R$ of the ball.

The third example, again in $\mathbb{R}^{3}$, investigates the case of (hyper)planes making small angles with each other. It is known that such case leads to
very slow convergence by using the existing POCS algorithm. We considered four planes through the $z$-axis and four planes through the $y$-axis, having the origin as intersection point. The equations of the planes in $\mathbb{R}^{3}$ were given by $y=x, y=1.4 x, y=1.7 x, y=2 x, z=4 x, z=4.4 x$, $z=4.7 x, z=5 x$ respectively.

| Starting Point | $(-3,0)$ | $(10,-10)$ | $(3,4)$ |
| :---: | :---: | :---: | :---: |
| PP | 1 | $3.279208 \times 10^{-3}$ | $3.661634 \times 10^{-3}$ |
|  |  | $5.000838 \times 10^{-4}$ | $5.49556 \times 10^{-4}$ |
| PAR | $9.972098 \times 10^{-3}$ | 4 | $1.129448 \times 10^{-2}$ |
| $3.128052 \times 10^{-3}$ |  | $3.427267 \times 10^{-3}$ |  |
| NEW0 | 9 | 5 | 9 |
| NEW3 | 20 | 5 | 8 |


| Starting Point | $(-17,12)$ | $(-2,1)$ | $(-100,-50)$ |
| :---: | :---: | :---: | :---: |
| PP | $3.601907 \times 10^{-3}$ | $3.202676 \times 10^{-3}$ | 1 |
|  | $5.419265 \times 10^{-4}$ | $4.89951 \times 10^{-4}$ |  |
| PAR | $1.185358 \times 10^{-2}$ | $9.768488 \times 10^{-3}$ | $8.859039 \times 10^{-3}$ |
|  | $3.548027 \times 10^{-3}$ | $3.080129 \times 10^{-3}$ | $2.859947 \times 10^{-3}$ |
| NEW0 | 10 | 10 | 10 |
| NEW3 | 8 | 9 | 42 |


| Starting Point | $(2,-4)$ | $(0,2)$ |
| :---: | :---: | :---: |
| PP | $3.005983 \times 10^{-3}$ | $3.694175 \times 10^{-3}$ |
|  | $4.637248 \times 10^{-4}$ | $5.537283 \times 10^{-4}$ |
| PAR | 5 | $9.757404 \times 10^{-3}$ |
|  |  | $3.077506 \times 10^{-3}$ |
| NEW0 | 5 | 9 |
| NEW3 | 5 | 10 |

Table 1.
In Table 2 we give, again for the four algorithms mentioned in the first example, the result either of the sum of the distances of the current iteration point to the eight planes after 1000 iterations, or the number of iterations needed to stop the procedure (this happened when the sum of the distances was less than $1 \times 10^{-8}$ ). Of course, this example is not
intended to present an alternative for existing algebraic methods in this case, but only to compare the behaviour of the mentioned algorithms.

| Starting Point | $(0.1,0.2,0.3)$ | $(-1,2,-3)$ | $(3,-1,2)$ |
| :---: | :---: | :---: | :---: |
| PP | $4.846649 \times 10^{-6}$ | $3.737408 \times 10^{-5}$ | $3.23111 \times 10^{-5}$ |
| PAR | $7.679005 \times 10^{-3}$ | $7.220158 \times 10^{-2}$ | $4.867536 \times 10^{-3}$ |
| NEW0 | $9.224179 \times 10^{-3}$ | 0.067403 | $6.159196 \times 10^{-2}$ |
| NEW3 | 374 | 372 | 430 |

Table 2.
The less good behaviour of NEW0 when compared to NEW3 may be explained by the fact that, with the relaxation used in NEW0, the monotoneous way of iteration has never been interrupted, contrary to what happens when using NEW3. In any case, this last example seems to suggest that interrupting the monotony by using suitable relaxations with respect to $\mathcal{D}$ can substantially improve the speed of convergence.

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