## Research Article

# A New Method for Proving Existence Theorems for Abstract Hammerstein Equations 

C. E. Chidume, ${ }^{1}$ C. O. Chidume, ${ }^{2}$ and Ma'aruf Shehu Minjibir ${ }^{1,3}$<br>${ }^{1}$ African University of Science and Technology, Abuja 900241, Nigeria<br>${ }^{2}$ Department of Mathematics and Statistics, Auburn University, Auburn, AL 36849-5168, USA<br>${ }^{3}$ Department of Mathematical Sciences, Bayero University, Kano 700241, Nigeria

Correspondence should be addressed to C. E. Chidume; cchidume@aust.edu.ng
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An abstract Hammerstein equation is an equation of the form $u+K F u=0$. A new method is introduced to prove the existence of a solution of this equation where $K$ and $F$ are nonlinear accretive (monotone) operators. The method does not involve the complicated technique of factorizing a linear map via a Hilbert space and does not involve the use of deep variational techniques.

## 1. General Introduction

Let $E$ be a real normed space and let $S:=\{x \in E:\|x\|=1\}$. The space $E$ is said to have Gâteaux differentiable norm if the limit

$$
\begin{equation*}
\lim _{t \rightarrow 0} \frac{\|x+t y\|-\|x\|}{t} \tag{1}
\end{equation*}
$$

exists for all $x, y \in S$; in this case $E$ is said to be smooth. $E$ is said to have uniformly Gâteaux differentiable norm if, for each $y \in S$, the limit is attained uniformly for $x \in S$. Further, $E$ is said to be uniformly smooth if the limit is attained uniformly for $(x, y) \in S \times S$. The modulus of smoothness of $E$, $\rho_{E}:[0, \infty) \rightarrow[0, \infty)$, is defined by

$$
\begin{aligned}
& \rho_{E}(\tau) \\
& \qquad:=\sup \left\{\frac{\|x+y\|+\|x-y\|}{2}-1:\|x\|=1,\|y\|=\tau\right\} ; \\
& \tau>0 .
\end{aligned}
$$

$E$ is equivalently said to be smooth if $\rho_{E}(\tau)>0 \forall \tau>0$. Let $q>1 ; E$ is said to be $q$-uniformly smooth (or to have a modulus
of smoothness of power type q) if there exists $c>0$ such that $\rho_{E}(\tau) \leq c \tau^{q}$.
$L_{p}, l_{p}$, and the Sobolev space $W_{m}^{p}, 1<p<\infty$, are all $q$-uniformly smooth. In fact

$$
L_{p} \text { or } l_{p} \text { or } W_{m}^{p} \text { is } \begin{cases}p \text {-uniformly smooth, } & 1<p \leq 2  \tag{3}\\ 2 \text {-uniformly smooth, } & p \geq 2\end{cases}
$$

Furthermore (see, e.g., [1]),

$$
\begin{align*}
\rho_{L_{p}} & (\tau) \\
& =\rho_{l_{p}}(\tau)=\rho_{W_{m}^{p}}(\tau) \\
& = \begin{cases}\left(1+\tau^{p}\right)^{1 / p}-1<\frac{1}{p} \tau^{p}, & 1<p \leq 2 \\
\frac{(p-1)}{2} \tau^{2}+o\left(\tau^{2}\right)<\frac{p-1}{2} \tau^{2}, & p \geq 2\end{cases} \tag{4}
\end{align*}
$$

Let $J_{q}$ denote the generalized duality mapping from $E$ to $2^{E^{*}}$ defined by

$$
\begin{equation*}
J_{q}(x):=\left\{f \in E^{*}:\langle x, f\rangle=\|x\|^{q},\|f\|=\|x\|^{q-1}\right\} \tag{5}
\end{equation*}
$$

where $E^{*}$ denotes the dual space of $E$ and $\langle\cdot, \cdot\rangle$ denotes the generalized duality pairing. It is well known (see, e.g., Xu [2]) that $J_{q}(x)=\|x\|^{q-2} J(x)$ if $x \neq 0$ where $J$ denotes $J_{2}$ (called the normalized duality mapping). It is well known that if $E^{*}$ is strictly convex, $J$ is single-valued. For more information and examples concerning (generalized) duality mappings, one may see the book of Cioranescu [3] and its review by Reich [4]. In the sequel, we will denote the singlevalued duality map by $j$.

A map $A: D(A) \subset X \rightarrow X$ is called accretive if, for all $x, y \in D(A)$, there exists $j(x-y) \in J(x-y)$ such that the following inequality holds:

$$
\begin{equation*}
\langle A x-A y, j(x-y)\rangle \geq 0 \tag{6}
\end{equation*}
$$

If $X$ is a real Hilbert space, the map $A$ is called monotone. In this case, $A$ satisfies the following condition:

$$
\begin{equation*}
\langle A x-A y, x-y\rangle \geq 0 \tag{7}
\end{equation*}
$$

The map $A$ is called strongly accretive if there exists $c>0$ such that, for all $x, y \in D(A)$, there exists $j(x-y) \in J(x-y)$, such that

$$
\begin{equation*}
\langle A x-A y, j(x-y)\rangle \geq c\|x-y\|^{2} \tag{8}
\end{equation*}
$$

A nonlinear integral equation of Hammerstein type (see, e.g., Hammerstein [5]) has the form

$$
\begin{equation*}
u(x)+\int_{\Omega} k(x, y) f(y, u(y)) d y=h(x) \tag{9}
\end{equation*}
$$

where $d y$ is a $\sigma$-finite measure on $\Omega$; the kernel $k$ is defined on $\Omega \times \Omega, f$ is a real-valued function defined on $\Omega \times \mathbb{R}$ and is, in general, nonlinear, and $h$ is a function on $\Omega$. Setting

$$
\begin{equation*}
K v(\cdot):=\int_{\Omega} k(\cdot, y) v(y) d y \quad \text { on } \Omega \tag{10}
\end{equation*}
$$

and $F u(\cdot):=f(\cdot, u(\cdot))$ on $\Omega$, then integral equation (9) can be put in abstract operator form as follows:

$$
\begin{equation*}
u+K F u=0 \tag{11}
\end{equation*}
$$

where, without loss of generality, we have taken $h \equiv 0$.
Interest in (9) stems mainly from the fact that several problems that arise in differential equations, for instance, elliptic boundary value problems whose linear parts possess Green's function, can, as a rule, be transformed into the form of (9).

Furthermore, equations of Hammerstein type play crucial role in the theory of optimal control systems, in automation, and in network theory (see, e.g., Dolezale [6]).

Several existence theorems for the solution of (9) have been proved by a host of distinguished mathematicians using various techniques (see, e.g., Browder and Gupta [7, 8], Chepanovich [9], and Petryshyn and Fitzpatrick [8]). In the remaining part of this section, we highlight the techniques used by Browder and Gupta [7] and Petryshyn and Fitzpatrick [8]. To do this, we first give definitions of some terms which are required in the theorems.

In the sequel, the symbol " $\rightarrow$ " denotes strong convergence while " $\rightarrow$ " denotes weak convergence.

Definition 1 (see, e.g., [7]). A mapping $A: D(A) \subset X^{*} \rightarrow X$ is said to be hemicontinuous if it is continuous from each line segment of $X^{*}$ to the weak topology of $X$. That is, $\forall u \in D(A)$, $\forall v \in X^{*}$, and $\left(t_{n}\right)_{n \geq 1} \subset \mathbb{R}^{+}$such that $t_{n} \rightarrow 0^{+}$and $u+t_{n} v \in$ $D(A)$ for $n$ sufficiently large and we have $A\left(u+t_{n} v\right) \rightharpoonup A(u)$.

Definition 2 (see, e.g., [7]). Let $A: X \rightarrow X^{*}$ be a bounded monotone linear mapping. $A$ is said to be anglebounded with constant $c \geq 0$ if, for all $u, v$ in $X, \mid\langle A u, v\rangle-$ $\langle A v, u\rangle \mid \leq 2 c\{\langle A u, u\rangle\}^{1 / 2}\{\langle A v, v\rangle\}^{1 / 2}$. (This is well defined since $\langle A u, u\rangle \geq 0$ and $\langle A v, v\rangle \geq 0$ by the linearity and monotonicity of $A$.)

In [7] Browder and Gupta proved the following theorem.
Theorem 3 (Browder-Gupta [7]). Let $X$ be a real Banach space and $X^{*}$ its conjugate dual space. Let $K$ be a monotone angle-bounded continuous linear mapping of $X$ into $X^{*}$ with constant of angle-boundedness $c \geq 0$. Let $F$ be a hemicontinuous (possibly nonlinear) mapping of $X^{*}$ into $X$ such that, for a given constant $k \geq 0$,

$$
\begin{equation*}
\left\langle v_{1}-v_{2}, F v_{1}-F v_{2}\right\rangle \geq-k\left\|v_{1}-v_{2}\right\|_{X^{*}}^{2} \tag{12}
\end{equation*}
$$

for all $v_{1}$ and $v_{2}$ in $X^{*}$. Suppose finally that there exists a constant $R$ with $k\left(1+c^{2}\right) R<1$ such that for $u$ in $X$

$$
\begin{equation*}
\langle K u, u\rangle \leq R\|u\|_{X}^{2} . \tag{13}
\end{equation*}
$$

Then, there exists exactly one solution $w$ in $X^{*}$ of the nonlinear equation

$$
\begin{equation*}
w+K F w=0 \tag{14}
\end{equation*}
$$

The main tool used by the authors in proving Theorem 3 is that of splitting the linear operator $K$ via a Hilbert space and then applying a deep result of Minty [10]. Precisely, they proved that if $X$ is a real Banach space, $X^{*}$ is its dual space, and $K$ is a bounded linear mapping of $X$ into $X^{*}$ which is monotone and angle-bounded, then there exist a Hilbert space $H$, a continuous linear mapping $S$ of $X$ into $H$ with adjoint $S^{*}$ injective, and a bounded skew-symmetric linear mapping $B$ of $H$ into $H$ such that

$$
\begin{equation*}
K=S^{*}(I+B) S \tag{15}
\end{equation*}
$$

(see Figure 1).
This factorization enabled the authors to transform the problem into another problem in a Hilbert space such that Hammerstein equation (11) has a solution if and only if the new problem has a solution in a real Hilbert space. They set $f=(I+B)^{-1}+K F K^{*}, D:=B(0,1)$, the closed unit ball in $H$, and showed that $f$ is hemicontinuous and monotone and satisfies $\langle u, f(u)\rangle \geq 0 \forall u \in D$. With these facts, they used the following result of Minty [10] to prove Theorem 3 (see [10] for definitions of terms).

Theorem 4 (Minty [10]). Let $D \subset X$ be bounded and surround 0 ; let $C \subset X$ contain $\overline{\operatorname{co}}(D)$ and surround every point of $\overline{\mathrm{co}}(D)$ densely. Let

$$
\begin{equation*}
f: C \longrightarrow X^{*} \tag{16}
\end{equation*}
$$


be monotone and hemicontinuous at every point of $\overline{\mathrm{co}}(D)$ and suppose

$$
\begin{equation*}
u \in D \text { implies }\langle u, f(u)\rangle \geq 0 \tag{17}
\end{equation*}
$$

Then, there exists $u \in \overline{\mathrm{co}}(D)$ such that $f(u)=0$.
Petryshyn and Fitzpatrick employed deep variational techniques to prove the existence of a solution to (11). They proved the following theorems.

Theorem 5 (Petryshyn-Fitzpatrick [8]). Let $X$ be a reflexive Banach space and let $K$ be a linear, monotone, and symmetric mapping of $X$ into $X^{*}$. Suppose $f$ is a weakly (sequential) lower semicontinuous functional on $X^{*}$ such that

$$
\begin{equation*}
f(u) \geq-\frac{1}{2} a_{1}\|u\|^{2}-a_{2}\|u\|^{\delta}-a_{3} \tag{18}
\end{equation*}
$$

where $a_{1}\|K\|<1, a_{2}>0, a_{3}>0$, and $0<\delta<2$. Suppose also that $F: X^{*} \rightarrow X$ is such that $\operatorname{grad}(f)=F$. Then,

$$
\begin{equation*}
w+K F w=0 \tag{19}
\end{equation*}
$$

has a solution in $X^{*}$.
Theorem 6 (Petryshyn-Fitzpatrick [8]). Let X be a reflexive Banach space with $K: X \rightarrow X^{*}$ linear, monotone, and symmetric. Let $F: X^{*} \rightarrow X$ be potential and have a Gâteaux derivative which satisfies the inequality

$$
\begin{equation*}
D F(u, v, v) \geq-a\|v\|^{2} \quad\left(v, u \in X^{*}\right) \tag{20}
\end{equation*}
$$

and $D N(t u, v, v)$ is continuous in $t \in[0,1]$ for $u$ and $v$ fixed, where $a\|K\|<1$. Then, (19) has a solution in $X^{*}$.

In this paper, we introduce a new method, perhaps simpler than methods used so far in the literature, of proving existence of solutions of Hammerstein equation in certain cases. To achieve this, we recast (11) into a fixed point problem and use a technique recently introduced by Chidume and Zegeye [11], some existence results of Deimling [12] for zeros of accretive maps, and some surjectivity results of Browder [13] for Lipschitz strongly accretive maps. No linearity assumption is imposed on any of our maps.

## 2. Preliminaries

Let $X$ be a normed linear space and let $K$ be a convex subset of $X$. For $x \in X$, the inward set, $I_{k}(x)$, of $x$ relative to $K$, is defined as follows:

$$
\begin{equation*}
I_{K}(x)=\{x+c(u-x): c \geq 1, u \in K\} \tag{21}
\end{equation*}
$$

A mapping $T: K \rightarrow X$ is said to be inward if $T x \in I_{K}(x)$ for each $x \in K$ and weakly inward if $T x$ belongs to the closure of $I_{K}(x)$ for each $x \in K$.

A relationship between the weak inward condition and the condition

$$
\begin{equation*}
\lim _{\lambda \rightarrow 0^{+}} \frac{\operatorname{dist}(x-\lambda A x, D(A))}{\lambda}=0 \quad \forall x \in D(A) \tag{22}
\end{equation*}
$$

for a map $A: D(A) \subset X \rightarrow X$ is given in Lemma 11. Further relationship between condition (22), the weak inward condition, and Lemma 11 can be found in [14].

In the sequel, $X$ is a $q$-uniformly smooth real Banach space, $q>1$, and $E:=X \times X$ with

$$
\begin{equation*}
\|[u, v]\|_{E}=\left(\|u\|^{q}+\|v\|^{q}\right)^{1 / q} \quad \forall[u, v] \in E . \tag{23}
\end{equation*}
$$

If $X(=H)$ is a real Hilbert space, we will denote $E$ by $E^{H}:=$ $H \times H$.

If $F$ and $K$ are maps from $X$ to $X$ such that range of $F$ is contained in domain of $K$, that is, $R(F) \subseteq D(K)$, Chidume and Zegeye [11] defined a map $A: E \rightarrow E$ as follows:

$$
\begin{equation*}
A[u, v]=[F u-v, K v+u] \tag{24}
\end{equation*}
$$

for all $u, v \in X$ and observed that $A[u, v]=0$ if and only if

$$
\begin{align*}
& F u-v=0 \\
& K v+u=0 \tag{25}
\end{align*}
$$

so that $u$ solves (11). System (25) can be recast as a fixed point problem as follows:

$$
\binom{u}{v}=\left(\begin{array}{cc}
0 & -K  \tag{26}\\
F & 0
\end{array}\right)\binom{u}{v} .
$$

We will use the ideas of map $A$ on $E$.
In Lemmas 9 and 10, we use the following variant definition of accretive maps as given by Deimling [12].

Definition 7 (accretive map in the sense of Deimling [12]). Let $X$ be real Banach space. A map $A: D(A) \subset X \rightarrow X$ is said to be accretive (in the sense of Deimling) if

$$
\begin{equation*}
\langle A(x)-A(y), x-y\rangle_{+} \geq 0 \quad \forall x, y \in D(A) \tag{27}
\end{equation*}
$$

where

$$
\begin{equation*}
\langle x, y\rangle_{+}:=\sup _{j(y) \in J(y)}\langle x, j(y)\rangle, \quad \forall x, y \in X . \tag{28}
\end{equation*}
$$

It is evident that, in any real Banach space, an accretive map is also accretive in the sense of Deimling. The converse is true in any real Banach $X$ whose dual $X^{*}$ is strictly convex or whose normalized duality map is single-valued. This is certainly the case when $X$ is $q$-uniformly smooth, $q>1$.

Definition 8 (see, e.g., [15]). A bounded convex subset $K$ of a Banach space $X$ is said to have normal structure if every convex subset $C$ of $K$ having more than one element contains at least one nondiametral point; that is, there exists $x^{0} \in C$ such that

$$
\begin{align*}
& \sup \left\{\left\|x^{0}-x\right\|: x \in C\right\}  \tag{29}\\
& \quad<\sup \{\|x-y\|: x, y \in C\}=d(C)
\end{align*}
$$

The Banach space $X$ is said to have normal structure if every bounded convex subset of $X$ has normal structure.

Lemma 9 (Deimling [12]). Let X be a reflexive real Banach space with normal structure and let $D$ be a closed convex bounded subset of $X$. Let $A: D \rightarrow X$ be a Lipschitz and accretive map satisfying condition (22). Then, $0 \in A(D)$.

Lemma 10 (Deimling [12]). Let $X$ be real Banach space and let $D$ be a closed convex subset of $X$. Let $A: D \subset X \rightarrow X$ be an accretive continuous map such that $\langle A x, x\rangle_{+} \geq 0$ for all $x \in X$ with $\|x\| \geq R$ for some $R>0$ or $\lim \|A x\|=\infty$ as $\|x\| \rightarrow \infty$. Suppose A satisfies condition (22) and suppose that $A(D)$ is closed. Then, $0 \in A(D)$.

Lemma 11 (Caristi [16]). Let D be a convex subset of a normed linear space $X$ and let $A: D \rightarrow X$ be a map. Then condition (22) holds if and only if $(I-A)$ is weakly inward and $I$ is the identity map on $D$.

Remark 12. In view of Lemma 11, if $D=H$ in Lemma 10, then condition (22) can be dropped.

Lemma 13 (Xu [2]). Let $q>1$ and $E$ a smooth real Banach space. Then the following are equivalent.
(i) E is q-uniformly smooth.
(ii) There exists a constant $d_{q}>0$ such that, for all $x, y \in E$,

$$
\begin{equation*}
\|x+y\|^{q} \leq\|x\|^{q}+q\left\langle y, j_{q}(x)\right\rangle+d_{q}\|y\|^{q} . \tag{30}
\end{equation*}
$$

(iii) There exists a constant $c_{q}>0$ such that for all $x, y \in E$ and $\lambda \in[0,1]$

$$
\begin{align*}
\|(1-\lambda) x+\lambda y\|^{q} \geq & (1-\lambda)\|x\|^{q}+\lambda\|y\|^{q} \\
& -w_{q}(\lambda) c_{q}\|x-y\|^{q}, \tag{31}
\end{align*}
$$

where $w_{q}(\lambda):=\lambda^{q}(1-\lambda)+\lambda(1-\lambda)^{q}$.
From now on, $c_{q}$ and $d_{q}$ denote the constants appearing in Lemma 13.

Lemma 14 (Chidume [15], p. 173). Let $X$ be a q-uniformly smooth real Banach space. Let $F, K: X \rightarrow X$ be maps with $F$ surjective such that the following conditions hold:
(i) there exists $\alpha>0$ such that, for each $u_{1}, u_{2} \in D(F)$,

$$
\begin{equation*}
\left\langle F u_{1}-F u_{2}, j_{q}\left(u_{1}-u_{2}\right)\right\rangle \geq \alpha\left\|u_{1}-u_{2}\right\|^{q} ; \tag{32}
\end{equation*}
$$

(ii) there exists $\beta>0$ such that, for each $u_{1}, u_{2} \in D(K)$,

$$
\begin{equation*}
\left\langle K u_{1}-K u_{2}, j_{q}\left(u_{1}-u_{2}\right)\right\rangle \geq \beta\left\|u_{1}-u_{2}\right\|^{q} ; \tag{33}
\end{equation*}
$$

(iii) $\left(1+d_{q}\right)\left(1+c_{q}\right) \geq 2^{q}, \min \{\alpha, \beta\}=: \gamma>\left(\left(1+d_{q}\right)(1+\right.$ $\left.\left.c_{q}\right)-2^{q}\right) / q\left(1+c_{q}\right)$.

Let a map A:E $\rightarrow$ E be defined by (24). Then, for each $z_{1}$, $z_{2} \in E$,

$$
\begin{align*}
& \left\langle A z_{1}-A z_{2}, j_{q}\left(z_{1}-z_{2}\right)\right\rangle \\
& \quad \geq\left[\gamma-q^{-1}\left(\left(1+d_{q}\right)-\frac{2^{q}}{\left(1+c_{q}\right)}\right)\right]\left\|z_{1}-z_{2}\right\|^{q} . \tag{34}
\end{align*}
$$

Lemma 15. Let $H$ be a real Hilbert space. Let $K: D(K) \subset$ $H \rightarrow H, F: D(F) \subset H: \rightarrow H$ be two monotone maps such that $R(F) \subset D(K)$. Then the map $A: D(F) \times D(K) \subset E^{H} \rightarrow$ $E^{H}$ defined by (24) is monotone.

Proof. The proof follows from the lines of argument of the proof of Lemma 14 (see Chidume and Zegeye [11]).

Lemma 16 (Chidume [15], p. 173). Let $X$ be a q-uniformly smooth real Banach space and let $K: D(K) \subset X \rightarrow X$, $F: D(F) \subset X \rightarrow X$ be two Lipschitz maps such that $R(F) \subset D(K)$. Let $A: D(A) \subset E$ be a map such that $D(F) \times D(K)=D(A)$ and defined by (24). Then, $A$ is Lipschitz.

We need the following definition which was given by Browder [17].

Definition 17 (Browder [17]). Let $X$ and $Y$ be real Banach spaces with $Y^{*}$ the conjugate space of $Y$. Let $\phi$ be a mapping of $X$ into $Y^{*}$ such that $\phi(X)$ is dense in $Y^{*}$ with

$$
\begin{align*}
\|\phi(x)\|_{Y^{*}} & =\|x\|  \tag{35}\\
\phi(\xi x) & =\xi \phi(x)
\end{align*}
$$

for all $x \in X, \xi \geq 0$. The mapping $f: X \rightarrow Y$ is said to be strongly $\phi$-accretive if there exists $c>0$ such that, for all $x$ and $u$ in $X$,

$$
\begin{equation*}
\langle f(x)-f(u), \phi(x-u)\rangle \geq c\|x-u\|^{2} \tag{36}
\end{equation*}
$$

It follows from this definition that if $X$ is a real Banach space such that the normalized duality map $J$ is single-valued and $J(X)$ is dense in $X^{*}$ (e.g., when $X$ is a reflexive and smooth real Banach space), then a strongly accretive map $A: X \rightarrow X$ is $J$-strongly accretive.

Theorem 18 (Browder [13]). Let $X$ and $Y$ be Banach spaces with $Y^{*}$ uniformly convex and suppose $f: X \rightarrow Y$ is a strongly $\phi$-accretive mapping satisfying a Lipschitz condition on each bounded subset of $X$. Then, $f(X)=Y$.

The following corollary follows from Theorem 18.

Corollary 19. Let $X$ be a real Banach space with uniformly convex dual $X^{*}$ and suppose $f: X \rightarrow X$ is a strongly accretive Lipschitz mapping. Then, $f(X)=X$.

## 3. Main Results

Let $X:=L_{p}, 1<p<2$, and let $E:=X \times X$ with $\|z\|_{E}^{2}:=$ $\|[u, v]\|_{E}^{2}=\|u\|_{X}^{2}+\|v\|_{E}^{2}$ for arbitrary $z=[u, v] \in E$. For $L_{p}$ spaces, $1<p<2$, the following estimate has been established (see, e.g., Chidume [15], p. 183):

$$
\begin{align*}
& A\left(u_{1}, u_{2}, v_{1}, v_{2}\right) \\
& \qquad:=\left[\left\langle v_{1}-v_{2}, j\left(u_{1}-u_{2}\right)\right\rangle+\left\langle u_{1}-u_{2}, j\left(u_{2}-u_{1}\right)\right\rangle\right] \\
& \quad \leq p(2-p)\left(\left\|u_{1}-u_{2}\right\|^{2}+\left\|v_{1}-v_{2}\right\|^{2}\right)  \tag{37}\\
& \forall u_{1}, u_{2}, v_{1}, v_{2} \in X .
\end{align*}
$$

We begin with a proof of the following theorem for $L_{p}$ spaces, $1<p<2$, which is new.

Theorem 20. Let $X=L_{p}(1<p<2)$; let $F, K: X \rightarrow X$ be mappings such that $D(K)=F(X)=X$ and the following conditions hold:
(a) there exists $\alpha>0$ such that, for each $u_{1}, u_{2} \in X$,

$$
\begin{equation*}
\left\langle F u_{1}-F u_{2}, j\left(u_{1}-u_{2}\right)\right\rangle \geq \alpha\left\|u_{1}-u_{2}\right\|^{2} \tag{38}
\end{equation*}
$$

(b) there exists $\beta>0$ such that, for each $u_{1}, u_{2} \in X$,

$$
\begin{equation*}
\left\langle K u_{1}-K u_{2}, j\left(u_{1}-u_{2}\right)\right\rangle \geq \beta\left\|u_{1}-u_{2}\right\|^{2} ; \tag{39}
\end{equation*}
$$

(c) $\gamma:=\min \{\alpha, \beta\}$ with $\gamma>p(2-p)$.

Let $E:=X \times X$ and define $A: E \rightarrow E$ by (24) for all $[u, v] \in E$. Then, for arbitrary $z_{1}, z_{2} \in E$, the following inequality holds:

$$
\begin{align*}
& \left\langle A z_{1}-A z_{2}, j^{E}\left(z_{1}-z_{2}\right)\right\rangle  \tag{40}\\
& \quad \geq[\gamma-p(2-p)]\left\|z_{1}-z_{2}\right\|^{2}
\end{align*}
$$

Proof. We compute as follows:

$$
\begin{aligned}
&\left\langle A z_{1}-A z_{2}, j^{E}\left(z_{1}-z_{2}\right)\right\rangle \\
&=\left\langle F u_{1}-F u_{2}, j\left(u_{1}-u_{2}\right)\right\rangle-\left\langle v_{1}-v_{2}, j\left(u_{1}-u_{2}\right)\right\rangle \\
&+\left\langle K v_{1}-K v_{2}, j\left(v_{1}-v_{2}\right)\right\rangle \\
&+\left\langle u_{1}-u_{2}, j\left(v_{1}-v_{2}\right)\right\rangle \\
& \geq \alpha\left\|u_{1}-u_{2}\right\|^{2}+\beta\left\|v_{1}-v_{2}\right\|^{2} \\
&-\left\langle v_{1}-v_{2}, j\left(u_{1}-u_{2}\right)\right\rangle+\left\langle u_{1}-u_{2}, j\left(v_{1}-v_{2}\right)\right\rangle \\
& \geq \gamma\left(\left\|u_{1}-u_{2}\right\|^{2}+\left\|v_{1}-v_{2}\right\|^{2}\right) \\
& \quad-\left[\left\langle v_{1}-v_{2}, j\left(u_{1}-u_{2}\right)\right\rangle-\left\langle u_{1}-u_{2}, j\left(v_{1}-v_{2}\right)\right\rangle\right]
\end{aligned}
$$

$$
\begin{align*}
\geq & \gamma\left\|z_{1}-z_{2}\right\|^{2}-A\left(u_{1}, u_{2}, v_{1}, v_{2}\right) \\
\geq & \gamma\left\|z_{1}-z_{2}\right\|^{2} \\
& \quad-p(2-p)\left(\left\|u_{1}-u_{2}\right\|^{2}+\left\|v_{1}-v_{2}\right\|^{2}\right) \\
= & (\gamma-p(2-p))\left\|z_{1}-z_{2}\right\|^{2} \quad \forall z_{1}, z_{2} \in E \tag{41}
\end{align*}
$$

completing proof of the theorem.
Remark 21. Observe that the condition $1+\sqrt{1-\gamma}<p<2$ implies $\gamma>p(2-p)$.

We now prove the following existence theorems.

### 3.1. The Case of Hilbert Spaces

Theorem 22. Let $H$ be a real Hilbert space and let $K$ : $D(K) \subset H \rightarrow H, F: D(F) \subset H \rightarrow H$ be two Lipschitz monotone maps such that $D(F)$ and $D(K)$ are closed, convex, and bounded and $R(F) \subset D(K)$. Let $A: D(A) \subset E^{H} \rightarrow$ $E^{H}$ be a map such that $D(F) \times D(K)=: D(A)$ and $A$ is defined by (24). Suppose that A satisfies condition (22). Then, Hammerstein equation (11) has a solution.

Proof. The fact that $K$ and $F$ are Lipschitz and monotone implies that $A$ is Lipschitz and monotone (Lemmas 15 and 16). Since the normalized duality map is the identity map in real Hilbert spaces, monotonicity of $A$ is equivalent to accretivity in the sense of Deimling. Also $D(A)$ is closed and convex since $D(F)$ and $D(K)$ are. Therefore, by Lemma 9, $0 \in A(D)$; that is, there exists $[u, v] \in D$ such that $F u-v=0$ and $K v+u=0$. So $u$ solves (11). This completes the proof.

Theorem 23. Let $H$ be a real Hilbert space and let $K: D(K) \subset$ $H \rightarrow H, F: D(F) \subset H \rightarrow H$ be two continuous monotone maps such that $D(F)$ and $D(K)$ are closed and convex and $R(F) \subset D(K)$. Let $A: D(A) \subset E^{H} \rightarrow E^{H}$ be a map such that $D(F) \times D(K)=: D(A)$ and $A$ is defined by (24). Suppose that $\langle A w, w\rangle \geq 0$ for all $w \in D(A)$ with $\|w\| \geq R$ for some $R>0$ or $\lim \|A w\|=\infty$ as $\|w\| \rightarrow \infty$ and suppose that $A$ satisfies condition (22). Suppose that $A(D(A))$ is closed. Then, Hammerstein equation (11) has a solution.

Proof. The fact that $K$ and $F$ are monotone implies that $A$ is monotone (Lemma 15). The fact that $D(F)$ and $D(K)$ are closed and convex implies that $D(A)$ is closed and convex. Also since $E^{H}$ is a real Hilbert space and the normalized duality map of any real Hilbert space is the identity map, we have $\langle A w, w\rangle_{+}=\langle A w, w\rangle$ for all $w \in D(A)$. Therefore, the assumptions on $A$ and $D(A)$ together with Lemma 10 give that $0 \in A(D)$; that is, there exists $[u, v] \in D$ such that $F u-v=0$ and $K v+u=0$. So $u$ solves (11). This completes the proof.

Corollary 24. Let $H$ be a real Hilbert space and let $K, F$ : $H \rightarrow H$ be two continuous monotone maps defined on $H$. Let $A: E^{H} \rightarrow E^{H}$ be a map defined by (24). Suppose that
$\langle A w, w\rangle \geq 0$ for all $w \in E^{H}$ with $\|w\| \geq R$ for some $R>0$ or $\lim \|A w\|=\infty$ as $\|w\| \rightarrow \infty$. Suppose that $A\left(E^{H}\right)$ is closed. Then, Hammerstein equation (11) has a solution.

Proof. Since $A$ is defined on $E^{H}$, it satisfies condition (22). Therefore, the result follows from Theorem 23.

### 3.2. The Case of $L^{p}$ Spaces, $1<p<\infty$

Theorem 25. Let $K: D(K) \subset L^{p} \rightarrow L^{p}$ and $F: D(F) \subset$ $L^{p} \rightarrow L^{p}$ be two Lipschitz mappings satisfying the following conditions:
(a) there exists $\alpha>0$ such that, for each $u_{1}, u_{2} \in D(F)$,

$$
\begin{equation*}
\left\langle F u_{1}-F u_{2}, j\left(u_{1}-u_{2}\right)\right\rangle \geq \alpha\left\|u_{1}-u_{2}\right\|^{2} ; \tag{42}
\end{equation*}
$$

(b) there exists $\beta>0$ such that, for each $u_{1}, u_{2} \in D(K)$,

$$
\begin{equation*}
\left\langle K u_{1}-K u_{2}, j\left(u_{1}-u_{2}\right)\right\rangle \geq \beta\left\|u_{1}-u_{2}\right\|^{2} \tag{43}
\end{equation*}
$$

Let $D(F)$ and $D(K)$ be closed, convex, and bounded such that $R(F) \subset D(K)$. Let $E:=L^{p} \times L^{p}$ and let $A: D(A) \subset E \rightarrow E$ be a map such that $D(F) \times D(K)=: D(A)$ and $A$ is defined by (24). Suppose that $A$ satisfies condition (22). Let $\gamma:=\min \{\alpha, \beta\}$. If $2 \leq p<\gamma+\sqrt{\gamma^{2}+4}$ or $1+\sqrt{1-\gamma}<p \leq 2$, then Hammerstein equation (11) has a solution.

Proof. The fact that $K$ and $F$ are Lipschitz implies that $A$ is Lipschitz by Lemma 16. Also $D(A)$ is closed and convex since $D(F)$ and $D(K)$ are.

Case $1\left(2 \leq p<\gamma+\sqrt{\gamma^{2}+4}\right)$. In this case $L^{p}$ is 2-uniformly smooth space and $c_{q}=d_{q}=p-1$ (see, e.g., [2]). Therefore, $\left(1+c_{q}\right)\left(1+d_{q}\right)=p^{2} \geq 4=2^{q}$ and

$$
\begin{equation*}
\gamma>\frac{1}{2 p}\left(p^{2}-4\right)=\frac{\left(1+d_{q}\right)\left(1+c_{q}\right)-2^{q}}{q\left(1+c_{q}\right)} \tag{44}
\end{equation*}
$$

for $2 \leq p<\gamma+\sqrt{\gamma^{2}+4}$. This implies by Lemma 14 that $A$ is accretive. Therefore, $A$ is accretive in the sense of Deimling. Hence, using Lemma 9, we have that $0 \in A(D)$; that is, there exists $[u, v] \in D$ such that $F u-v=0$ and $K v+u=0$. So $u$ solves (11).

Case $2(1+\sqrt{1-\gamma}<p \leq 2)$. The condition $1+\sqrt{1-\gamma}<$ $p \leq 2$ implies that $\gamma>p(2-p)$. Hence, by Theorem 20, $A$ is accretive. We conclude as in Case 1. This completes the proof.

Theorem 26. Let $K: D(K) \subset L^{p} \rightarrow L^{p}, F: D(F) \subset$ $L^{p} \rightarrow L^{p}$ be two continuous mappings satisfying the following conditions:
(a) there exists $\alpha>0$ such that, for each $u_{1}, u_{2} \in D(F)$,

$$
\begin{equation*}
\left\langle F u_{1}-F u_{2}, j\left(u_{1}-u_{2}\right)\right\rangle \geq \alpha\left\|u_{1}-u_{2}\right\|^{2} ; \tag{45}
\end{equation*}
$$

(b) there exists $\beta>0$ such that, for each $u_{1}, u_{2} \in D(K)$,

$$
\begin{equation*}
\left\langle K u_{1}-K u_{2}, j\left(u_{1}-u_{2}\right)\right\rangle \geq \beta\left\|u_{1}-u_{2}\right\|^{2} \tag{46}
\end{equation*}
$$

Let $D(F)$ and $D(K)$ be closed and convex, such that $R(F) \subset$ $D(K)$. Let $E:=L^{p} \times L^{p}$ and let $A: D(A) \subset E \rightarrow E$ be a mapping such that $D(F) \times D(K)=: D(A)$ and $A$ is defined by (24) for $[u, v] \in D(A)$. Suppose that $\langle A w, w\rangle_{+} \geq 0$ for all $w \in D(A)$ with $\|w\| \geq R$ for some $R>0$ or $\lim \|A w\|=\infty$ as $\|w\| \rightarrow \infty$ and suppose A satisfies condition (22). Let $\gamma:=$ $\min \{\alpha, \beta\}$. If $2 \leq p<\gamma+\sqrt{\gamma^{2}+4}$ or $1+\sqrt{1-\gamma}<p \leq 2$, then Hammerstein equation (11) has a solution.

Proof. Evidently, continuity of $K$ and $F$ gives the continuity of $A$. Also $D(A)$ is closed and convex since $D(F)$ and $D(K)$ are. The rest follows as in the proof of Theorem 25. This completes the proof.

Corollary 27. Let $K: D(K) \subset L^{p} \rightarrow L^{p}, F: D(F) \subset$ $L^{p} \rightarrow L^{p}$ be two continuous accretive mappings satisfying the following conditions:
(a) there exists $\alpha>0$ such that, for each $u_{1}, u_{2} \in D(F)$,

$$
\begin{equation*}
\left\langle F u_{1}-F u_{2}, j\left(u_{1}-u_{2}\right)\right\rangle \geq \alpha\left\|u_{1}-u_{2}\right\|^{2} \tag{47}
\end{equation*}
$$

(b) there exists $\beta>0$ such that, for each $u_{1}, u_{2} \in D(K)$,

$$
\begin{equation*}
\left\langle K u_{1}-K u_{2}, j\left(u_{1}-u_{2}\right)\right\rangle \geq \beta\left\|u_{1}-u_{2}\right\|^{2} \tag{48}
\end{equation*}
$$

Let $D(F)=L^{p}=D(K)$. Let $E:=L^{p} \times L^{p}$ and let $A: E \rightarrow$ $E$ be a mapping defined by (24) $[u, v] \in D(A)$. Suppose that $\langle A w, w\rangle_{+} \geq 0$ for all $w \in E$ with $\|w\| \geq R$ for some $R>0$ or $\lim \|A w\|=\infty$ as $\|w\| \rightarrow \infty$. Let $\gamma:=\min \{\alpha, \beta\}$. If $2 \leq$ $p<\gamma+\sqrt{\gamma^{2}+4}$ or $1+\sqrt{1-\gamma}<p \leq 2$, then Hammerstein equation (11) has a solution.

Proof. Since $A$ is defined on $E$, it satisfies condition (22) of Theorem 26. Also $D(A)$ is closed and convex. Therefore, the result follows from Theorem 26.

### 3.3. The Case of Hilbert Spaces with Lipschitz Strongly Monotone Mappings

Theorem 28. Let $H$ be a real Hilbert space and let $K: H \rightarrow$ $H, F: H \rightarrow H$ be two Lipschitz strongly monotone mappings with constants $\alpha, \beta$, respectively. Let $A: E^{H} \rightarrow E^{H}$ be a mapping defined by (24) for $[u, v] \in E^{H}$. Then, Hammerstein equation (11) has a solution.

Proof. Using Lemma 16 we have that $A$ is Lipschitz. Also since every real Hilbert space is $q$-uniformly smooth with $q=2$, $d_{q}=c_{q}=1$, we have that $\left(1+c_{q}\right)\left(1+d_{q}\right)=4=2^{q}$. Also $\min \{\alpha, \beta\}>0=\left(\left(1+c_{q}\right)\left(1+d_{q}\right)-2^{q}\right) / q$. Therefore, $A$ is strongly monotone by Lemma 14. Since $E^{H}$ is a real Hilbert space and every real Hilbert space is uniformly convex, we invoke Corollary 19 to obtain that $A\left(E^{H}\right)=E^{H}$. So there
exists $[u, v] \in E^{H}$ such that $A[u, v]=0$; that is, $F u-v=$ $0, K v+u=0$. Hence $u$ solves (11). This completes the proof.
3.4. The Case of $L_{p}$ Spaces, $1<p<\infty$, with Lipschitz Strongly Accretive Mappings

Theorem 29. Let $K: L^{p} \rightarrow L_{p}, F: L^{p} \rightarrow L_{p}$ be two Lipschitz mappings satisfying the following conditions:
(a) there exists $\alpha>0$ such that, for each $u_{1}, u_{2} \in L^{p}$,

$$
\begin{equation*}
\left\langle F u_{1}-F u_{2}, j\left(u_{1}-u_{2}\right)\right\rangle \geq \alpha\left\|u_{1}-u_{2}\right\|^{2} \tag{49}
\end{equation*}
$$

(b) there exists $\beta>0$ such that, for each $u_{1}, u_{2} \in L^{p}$,

$$
\begin{equation*}
\left\langle K u_{1}-K u_{2}, j\left(u_{1}-u_{2}\right)\right\rangle \geq \beta\left\|u_{1}-u_{2}\right\|^{2} \tag{50}
\end{equation*}
$$

Let $E:=L^{p} \times L^{p}$ and let $A: E \rightarrow E$ be a mapping defined by (24). Let $\gamma:=\min \{\alpha, \beta\}$. If $2 \leq p<\gamma+\sqrt{\gamma^{2}+4}$ or $1+\sqrt{1-\gamma}<$ $p \leq 2$, then, Hammerstein equation (11) has a solution.

Proof. Using Lemma 16 we have that $A$ is Lipschitz.
Case $1\left(2 \leq p<\gamma+\sqrt{\gamma^{2}+4}\right)$. In this case $L^{p}$ is 2-uniformly smooth space and $c_{q}=d_{q}=p-1$ (see, e.g., [2]). Therefore, $\left(1+c_{q}\right)\left(1+d_{q}\right)=p^{2} \geq 4=2^{q}$ and

$$
\begin{equation*}
\gamma>\frac{1}{2 p}\left(p^{2}-4\right)=\frac{\left(1+d_{q}\right)\left(1+c_{q}\right)-2^{q}}{q\left(1+c_{q}\right)} \tag{51}
\end{equation*}
$$

for $2 \leq p<\gamma+\sqrt{\gamma^{2}+4}$. This implies by Lemma 14 that $A$ is strongly accretive. Since every $L_{p}$ space, $2 \leq p<\gamma+\sqrt{\gamma^{2}+4}$, is uniformly convex, by Corollary $19, A\left(L_{p}\right)=L_{p}$. Therefore there exists $[u, v] \in D$ such that $A[u, v]=0$; that is, $F u-v=0$ and $K v+u=0$. So $u$ solves (11).

Case $2(1+\sqrt{1-\gamma}<p \leq 2)$. The inequality $1+\sqrt{1-\gamma}<$ $p \leq 2$ implies that $\gamma>p(2-p)$. Hence by Theorem $20 A$ is strongly accretive. The result now follows as in Case 1 since every $L_{p}$ space, $1+\sqrt{1-\gamma}<p \leq 2$, is uniformly convex. This completes the proof.

## Conflict of Interests

The authors declare that there is no conflict of interests regarding the publication of this paper.

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