A New Fixed Point Result and its Application to Existence Theorem for Nonconvex Hammerstein Type Integral Inclusions^{*}

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Abstract. In this paper, a generalization of Nadler's fixed point theorem is presented for H^+ -type k-multivalued weak contractive mappings. We consider a nonconvex Hammerstein type integral inclusion and prove an existence theorem by using an H^+ -type multi-valued weak contractive mapping.

Keywords and Phrases: Multi-valued contraction map, multi-valued weak contractive map, H^+ -type multi-valued weak contractive map, Hammerstein type integral inclusion, Fixed point. 2000 Mathematical Subject Classification : 47G20, 47H10, 47H20, 54H15, 54H25, 81Q05.

1. Introduction

In 1969, Nadler [16] proved a fixed point theorem for the set-valued contractions, which is of fundamental importance in nonlinear analysis. Inspired from the fixed point result of Nadler [16], the fixed point theory of set-valued contraction was further developed in different directions by many authors, in particular, by Reich [20, 21], Mizoguchi and Takahashi [15], Ciric [3], Kaneko [9], Lim [13], Lami Dozo [14], Feng and Liu [5], Klim and Wardowski [10], Suzuki [22], Pathak and Shahzad [17, 18] and many others. For details, see [19]. An interesting application of a consequence of Nadler's fixed point theorem was given in Cernea [2]. For other applications of the same result see, for example, [4] [6], [7], [8], [12] and [19].

2. Preliminaries and Definitions

Let (X, d) be a metric space. Let CB(X) and C(X) denote the collection of all nonempty closed and bounded subsets of X and the collection of all compact subsets of X, respectively.

For $A, B \in CB(X)$, let

$$H(A,B) = \max\Big\{\rho(A,B),\rho(B,A)\Big\},$$

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$$H^{+}(A,B) = \frac{1}{2} \Big\{ \rho(A,B) + \rho(B,A) \Big\},\$$

where $\rho(A, B) = \sup_{x \in A} d(x, B)$ and $d(x, B) = \inf_{y \in B} d(x, y)$. It is well known that H is a metric on CB(X). Such a map H is called *Pompeiu-Hausdorff metric* induced by d.

A mapping $T: X \to CB(X)$ is said to be a

• multi-valued contraction mapping if there exists a fixed real number k, 0 < k < 1 such that

$$H(Tx, Ty) \le k \, d(x, y), \tag{2.1}$$

for all $x, y \in X$.

• multi-valued weak contractive mapping if there exists a fixed real number k, 0 < k < 1 such that

 $H(Tx,Ty) \le k \, \max\{d(x,y), d(x,Tx), d(y,Ty), [d(x,Ty) + d(y,Tx)]/2\}, \tag{2.2}$ for all $x,y \in X.$

• multi-valued quasi-contraction mapping if there exists a fixed real number k, 0 < k < 1 such that

$$H(Tx, Ty) \le k \max\{d(x, y), d(x, Tx), d(y, Ty), d(x, Ty), d(y, Tx)\},$$
(2.3)

for all $x, y \in X$.

Proposition 2.1([18]). H^+ is a metric on CB(X).

Notice that the two metrics H and H^+ are equivalent [11] since

$$\frac{1}{2}H(A,B) \le H^+(A,B) \le H(A,B).$$

In the light of this equivalence and referring to Kuratowski [11], we conclude that $(CB(X), H^+)$ is complete whenever (X, d) is complete. Indeed, it is a simple consequence of the completeness of the Hausdorff metric H. Moreover, C(X) is a closed subspace of $(CB(X), H^+)$.

Notice also that $H^+ : \mathcal{CB}(X) \times \mathcal{CB}(X) \to \mathbf{R}$ is a continuous function. To see this, we observe that the inequality

$$H^+(A,B) \le H^+(A,C) + H^+(C,B)$$

holds for any $A, B, C \in \mathcal{CB}(X)$. Now pick any $(A_0, B_0) \in \mathcal{CB}(X) \times \mathcal{CB}(X)$. Then for a given $\epsilon > 0$, we can choose a positive number $\delta = \frac{\epsilon}{2}$ such that

$$|H^+(A,B) - H^+(A_0,B_0)| \le H^+(A,A_0) + H^+(B_0,B) < \delta + \delta = 2\delta = \epsilon$$

whenever $H^+(A, A_0) < \delta, H^+(B_0, B) < \delta$. This shows that H^+ is continuous at (A_0, B_0) .

In [16], S. B. Nadler proved the following result, which he announced earlier.

Theorem 2.2. Let (X, d) be a complete metric space and $T : X \to CB(X)$ a multi-valued contraction mapping. Then T has a fixed point.

In this paper, we intend to generalize this result by weakening the multi-valued contraction to an H^+ -type multi-valued weak contractive mapping. Our main result is summarized in Section 3. In Section 4, we consider a nonconvex Hammerstein type integral inclusion and prove an existence theorem by using an H^+ -type multi-valued weak contractive mapping.

3. Main results

We begin our discussion with the following definition.

Definition 3.1. Let (X, d) be a metric space. A multi-valued mapping $T : X \to \mathcal{CB}(X)$ is called H^+ -contraction if

(1) there exists a fixed real number k, 0 < k < 1 such that

 $H^+(Tx, Ty) \le kd(x, y)$ for every $x, y \in X$,

(2) for every x in X, y in T(x) and $\epsilon > 0$, there exists z in T(y) such that

$$d(y,z) \le H^+(T(y),T(x)) + \epsilon.$$

In [18], Pathak and Shahzad proved the following result.

Theorem 3.2. Every H^+ -type multi-valued contraction mapping $T : X \to CB(X)$ with Lipschitz constant 0 < k < 1 has a fixed point.

We now introduce the following definition.

Definition 3.3. Let (X, d) be a metric space. A mapping $T : X \to CB(X)$ is called an H^+ -type multi-valued weak contractive mapping if the condition (2) holds and there exists a fixed real number k, 0 < k < 1 such that

$$H^{+}(Tx, Ty) \le k \max\{d(x, y), d(x, Tx), d(y, Ty), [d(x, Ty) + d(y, Tx)]/2\},$$
(3.1)

for all x, y in X.

Now we state and prove our main result.

Theorem 3.4. Let (X, d) be a complete metric space and $T : X \to CB(X)$ an H^+ -type multivalued weak k-contractive mapping with 0 < k < 1. Then T has a fixed point.

Proof. Notice first that for each $A, B \in CB(X), a \in A$ and $\alpha > 0$ with $H^+(A, B) < \alpha$, there exists $b \in B$ such that $\max\{d(a, b), d(a, Ta), d(b, Tb), \frac{1}{2}[d(a, Tb) + d(b, Ta)]\} < \alpha$. Now, let L > 0 be such that k < L < 1. Then

$$H^{+}(Tx, Ty) < L \max\{d(x, y), d(x, Tx), d(y, Ty), [d(x, Ty) + d(y, Tx)]/2\},$$
(3.2)

for any $x, y \in X, x \neq y$.

Now we choose a sequence $\{x_n\}$ recursively in X in the following way. Let $x_0 \in X$ be arbitrary. Fix an element x_1 in Tx_0 . From (2) it follows that we can choose $x_2 \in Tx_1$ such that

$$d(x_1, x_2) \le H^+(Tx_0, Tx_1) + \epsilon \tag{3.3}$$

In general, if x_n be chosen, then we choose $x_{n+1} \in Tx_n$ such that

$$d(x_n, x_{n+1}) \le H^+(Tx_{n-1}, Tx_n) + \epsilon.$$
(3.4)

Set $\epsilon = (\frac{1}{\sqrt{L}} - 1)H^+(Tx_{n-1}, Tx_n)$. Then from (3.4), it follows that

$$d(x_n, x_{n+1}) \le H^+(Tx_{n-1}, Tx_n) + \left(\frac{1}{\sqrt{L}} - 1\right) H^+(Tx_{n-1}, Tx_n) = \frac{1}{\sqrt{L}} H^+(Tx_{n-1}, Tx_n).$$

Thus, we have

$$\sqrt{L} d(x_n, x_{n+1}) \le H^+(Tx_{n-1}, Tx_n)$$
(3.5)

for each $n \in \mathbf{N}$.

Thus, from (3.2) we have

$$\begin{split} \sqrt{L} \, d(x_n, x_{n+1}) &< L \, \max\{d(x_{n-1}, x_n), d(x_{n-1}, Tx_{n-1}), d(x_n, Tx_n), \\ & [d(x_{n-1}, Tx_n) + d(x_n, Tx_{n-1})]/2\} \\ &\leq (\sqrt{L})^2 \, \max\{d(x_{n-1}, x_n), d(x_{n-1}, x_n), d(x_n, x_{n+1}), d(x_{n-1}, x_{n+1})/2\} \\ &\leq (\sqrt{L})^2 \, \max\{d(x_n, x_{n-1}), d(x_n, x_{n+1}), [d(x_{n-1}, x_n) + d(x_n, x_{n+1})]/2\} \\ &= (\sqrt{L})^2 \, \max\{d(x_n, x_{n-1}), d(x_n, x_{n+1})\}. \end{split}$$

It follows that

$$d(x_n, x_{n+1}) < \sqrt{L} \max\{d(x_n, x_{n-1}), d(x_n, x_{n+1})\}$$
(3.6)

for each $n \in \mathbf{N}$. Note that if $x_n = x_{n+1}$ for some $n \in \mathbf{N}$, then $x_n = x_{n+1} \in Tx_n$, that is, x_n is a fixed point of T and we are finished. So, we may assume that $d(x_{n+1}, x_n) > 0$ for each $n \in \mathbf{N}$. Suppose that $d(x_{n-1}, x_n) < d(x_n, x_{n+1})$ for some $n \in \mathbf{N}$, then inequality (3.6) gives

$$d(x_n, x_{n+1}) < \sqrt{L} d(x_n, x_{n+1}),$$

a contradiction. So we must have $d(x_{n-1}, x_n) \ge d(x_n, x_{n+1})$ for each $n \in \mathbb{N}$. Hence, for all $n \in \mathbb{N}$, (3.6) yields

$$d(x_n, x_{n+1}) < c \, d(x_{n-1}, x_n), \tag{3.7}$$

where $c = \sqrt{L}$. Repeating the same argument n-times as in (3.7), we obtain

$$d(x_n, x_{n+1}) < c^n \, d(x_0, x_1). \tag{3.8}$$

It is obvious that $\{x_n\}$ is bounded. Indeed, for any $n \in \mathbf{N}$, we have

$$d(x_0, x_n) \le \sum_{i=0}^{n-1} d(x_i, x_{i+1}) < (1 + c + c^2 + \dots + c^n) d(x_0, x_1)$$
$$< (1 + c + c^2 + \dots) d(x_0, x_1) = \frac{1}{1 - c} d(x_0, x_1) < \infty$$

Further, by virtue of (3.8), one may observe that $\{x_n\}$ is a Cauchy sequence. Since X is complete, there exists $u \in X$ such that $\lim_{n\to\infty} x_n = u$. Assume that $u \notin Tu$, that is, d(u,Tu) > 0. Now using (3.2) we have

$$\frac{1}{2} \Big\{ \rho(Tx_n, Tu) + \rho(Tu, Tx_n) \Big\} = H^+(Tx_n, Tu)
< L \max\{d(x_n, u), d(x_n, Tx_n), d(u, Tu), [d(x_n, Tu) + d(u, Tx_n)]/2\}
\leq L \max\{d(x_n, u), d(x_n, x_{n+1}), d(u, Tu), [d(x_n, Tu) + d(u, x_{n+1})]/2\},$$

it follows that

$$\frac{1}{2}\liminf_{n\to\infty}\left\{\rho(Tx_n,Tu)+\rho(Tu,Tx_n)\right\}\leq L\,d(u,Tu).$$

Since $\lim_{n\to\infty} d(x_{n+1}, u) = 0$ exists, and

$$d(u, Tu) = \frac{1}{2}[d(u, Tu) + d(Tu, u)] \le \frac{1}{2}[\rho(Tx_n, Tu) + \rho(Tu, Tx_n)] + d(x_{n+1}, u),$$

it follows that

$$d(u,Tu) \leq \frac{1}{2} \liminf_{n \to \infty} [\rho(Tx_n,Tu) + \rho(Tu,Tx_n)] + \liminf_{n \to \infty} d(x_{n+1},u)$$

$$\leq L d(u,Tu) + \lim_{n \to \infty} d(x_{n+1},u) = L d(u,Tu) < d(u,Tu),$$

a contradiction. This implies that d(u, Tu) = 0, and, since Tu is closed, it must be the case that $u \in Tu$.

Notice that every multi-valued contraction mapping with respect to Pompeiu-Hausdorff metric H is an H^+ -type multi-valued weak contractive mapping but the converse implication need not be true. To see this, we have the following example:

Example 3.5. Let X = [-2, 2] and $d : X \times X \to \mathbf{R}$ be a standard metric. Let $T : X \to CB(X)$ be defined by $Tx = \{\frac{x}{4}\}$, if $x \in [-1, 2]$ and $Tx = \{2\}$, otherwise. It is clear that if $x, y \in [-1, 2]$ or $x, y \in [-2, -1)$, then

$$H^+(Tx,Ty) \le \frac{1}{4}d(x,y).$$

If $x \in [-1, 2]$ and $y \in [-2, -1)$, then we have

$$H^{+}(Tx,Ty) = \frac{1}{2}[|2 - \frac{x}{4}| + |2 - \frac{x}{4}|] = |2 - \frac{x}{4}| \le 2 + \frac{1}{4} = \frac{3}{4} \cdot 3 \le \frac{3}{4} \cdot \max\{d(y,Ty), d(x,Tx)\}.$$

It follows that

$$H^{+}(Tx, Ty) \le k \max\{d(x, y), d(x, Tx), d(y, Ty), [d(x, Ty) + d(y, Tx)]/2\}$$

for all $x, y \in X$ and $k \in [\frac{3}{4}, 1)$. To check the condition (2), we consider the following cases: Case 1. If $x \in [-2, -1)$, then for any $y \in Tx = \{2\}$, there exists $z \in Ty = \{\frac{1}{2}\}$ such that for any $\epsilon > 0$

$$d(y,z) = \frac{3}{2} \le \frac{3}{2} + \epsilon = H^+(Ty,Tx) + \epsilon.$$

Case 2. If $x \in [-1, 2]$, then for any $y \in Tx = \{\frac{x}{4}\}$, there exists $z \in Ty = \{\frac{x}{16}\}$ such that for any $\epsilon > 0$

$$d(y,z) = \frac{3|x|}{16} \le \frac{3|x|}{16} + \epsilon = H^+(Ty,Tx) + \epsilon.$$

Thus all the conditions of Theorem 3.4 are satisfied. Moreover, $0 \in T0 = \{0\}$ is a fixed point of T.

Notice that the map T does not satisfy the assumptions of Theorem 2.2 and Theorem 3.2. Indeed, for x = -1 and $y \to -1$ from the left we have

$$H(T(-1), T(y)) = H^+(T(-1), T(y)) = 2 + \frac{1}{4} > k \, d(-1, y),$$

for all $k \in (0, 1)$.

We also notice that since

$$[d(x, Ty) + d(y, Tx)]/2 \le \max\{d(x, Ty), d(y, Tx)\}$$

for all $x, y \in X$, it follows that every weak contractive mapping is quasi-contraction.

Using the technique of the proof of Theorem 3.4, one can easily prove the following result.

Theorem 3.6. Let (X, d) be a complete metric space. Let $T : X \to CB(X)$ be a H^+ -type k-multi-valued quasi-contraction mapping with $0 < k < \frac{1}{2}$. Then, T has a fixed point.

Pathak and Shahzad [18] introduced the class of H^+ -type nonexpansive mappings

Definition 3.7. Let $(X, \|\cdot\|)$ be a Banach space. A multi-valued map $T: X \to \mathcal{CB}(X)$ is called H^+ -nonexpansive if

(1')
$$H^+(Tx, Ty) \le ||x - y|| \text{ for every } x, y \in X,$$

(2') for every x in X, y in T(x) and $\epsilon > 0$, there exists z in T(y) such that

$$||y - z|| \le H^+(T(y), T(x)) + \epsilon.$$

Applying the main result of this section, we obtain the following result which plays a role in the next section.

Proposition 3.8.([18]). Let (X, d) be a complete metric space. Suppose that $T_i : X \to CB(X), i = 1, 2$, are two H^+ -type multi-valued contraction mappings with Lipschitz constant L < 1. Then if $Fix(T_1)$ and $Fix(T_2)$ denote the respective fixed point sets of T_1 and T_2 ,

$$H^+(Fix(T_1), Fix(T_2)) \le \frac{1}{1 - \sqrt{L}} \sup_{x \in X} H^+(T_1x, T_2y).$$

4. Existence Theorem for Nonconvex Hammerstein Type Integral Inclusions

Let $0 < T < \infty$, I := [0, T] and $\mathcal{L}(I)$ denote the σ -algebra of all Lebesgue measurable subsets of I. Let E be a real separable Banach space with the norm $\|\cdot\|$. Let $\mathcal{P}(E)$ denote the family of all nonempty subsets of E and $\mathcal{B}(E)$ the family of all Borel subsets of E.

In what follows, as usual, we denote by C(I, E) the Banach space of all continuous functions $x(\cdot) : I \to E$ endowed with the norm $||x(\cdot)||_C = \sup_{t \in I} ||x(t)||$. Consider the following integral equation

$$x(t) = \lambda(t) + \int_0^T k(t,s) g(t,s,u(s)) \, ds \text{ on } [0,T].$$
(4.1)

Here λ, k and g are given functions, where $\lambda(\cdot) : I \to E$ is a function with Banach space value, $k: I \times I \to \mathbf{R}_+ = [0, \infty)$ is a positive real single-valued function, while $g: I \times I \times E \to E$ is a map. Let $p \in [1, \infty), q \in [1, \infty)$, and let $r \in [1, \infty)$ be the conjugate exponent of q, that is 1/q + 1/r = 1. Let $\|\cdot\|_p$ denote the p-norm of the space $L^p(I, E)$ and is defined by $\|u\|_p = (\int_0^T \|u(s)\|^p ds)^{1/p}$ for all $u \in L^p(I, E)$. Consider the Nemitsky operator associated to g, p, q and $G: L^p(I, E) \to L^q(I, E)$ given by

$$G(u) = g(t, s, u(s)) a.e.$$
on I

Consider the linear integral operator of kernel $k, S: L^q(I, E) \to L^p(I, E)$ given by

$$S(u) = \lambda(t) + \int_0^T k(t,s)u(s)ds \ a.e. \text{ on } I.$$

Thus the Hammerstein type integral equation (4.1) is transformed into the form

$$x = SG(u), \quad u \in L^p(I, E) \text{ a.e. on } I \tag{4.1'}$$

$$u(t) \in F(t, V(x)(t))$$
 a.e. $(I := [0, T]),$ (4.2)

where $V : C(I, E) \to C(I, E)$ is a given mapping. In the sequel, we also use the following: For any $x \in E$, $\lambda \in C(I, E)$, $\sigma \in L^p(I, E)$, we define the set-valued maps $M_{\lambda,\sigma}(t) := F(t, V(x_{\sigma,\lambda})(t))$, $t \in I, T_{\lambda}(\sigma) := \{\psi(\cdot) \in L^p(I, E) : \psi(t) \in M_{\lambda,\sigma}(t) \ a.e. (I)\}.$

In order to study problem (4.1)-(4.2) we introduce the following assumption.

Hypothesis 4.1. Let $F(\cdot, \cdot) : I \times E \to \mathcal{P}(E)$ be a set-valued map with nonempty closed values satisfying:

(H₁) The function $k: I \times I \to \mathbf{R}_+$ satisfies that $k(t, \cdot) \in L^r(I)$, and $t \to ||k(t, \cdot)||_r \in L^p(I)$.

 (H_2) The set-valued map $F(\cdot, \cdot)$ is $\mathcal{L}(I) \otimes \mathcal{B}(E)$ measurable.

 (H_3) There exists $L(\cdot) \in L^1(I, \mathbf{R}_+)$ such that, for almost all $t \in I, F(t, \cdot)$ is L(t)-Lipschitz in the sense that

$$H^+(F(t,x), F(t,y)) \le L(t) ||x-y||$$
 (C1)

for all x, y in E, and for any $x, y \in X$, $w \in F(t, x)$ and any $\epsilon > 0$, there exists $z \in F(t, y)$ such that

$$||w - z||^p \le H^+(F(t, x), F(t, y)) + \epsilon$$
 (C2)

and $T_{\lambda}(\cdot)$ satisfies the condition: For any $\sigma \in L^{p}(I, E)$, $\sigma_{1} \in T_{\lambda}(\sigma)$ and any given $\epsilon > 0$, there exists $\sigma_{2} \in T_{\lambda}(\sigma_{1})$ such that

$$\|\sigma_1 - \sigma_2\|_p \le H^+(T_\lambda(\sigma), T_\lambda(\sigma_1)) + \epsilon.$$
(C3)

(*H*₄) The mappings $k: I \times I \to \mathbf{R}_+, g: I \times I \times E \to E$ are continuous, $V: C(I, E) \to C(I, E)$ and there exist constants $M_1, M_2, M_3 > 0$ such that

$$||g(t,s,u_1) - g(t,s,u_2)|| \le M_1 ||u_1 - u_2||^p, \ \forall u_1, u_2 \in E,$$
$$||V(x_1)(t) - V(x_2)(t)|| \le M_2 ||x_1(t) - x_2(t)||, \ \forall t \in I, \forall x_1, x_2 \in C(I,E),$$
$$|k(t,s)| \le M_3 \ \forall t, s \in I.$$

and

It is worth mentioning that the system (4.1)-(4.2) includes a large variety of differential inclusions and control systems.

Assume that U is an open bounded subset of \mathbf{R}^n (or Y, a subset of E homeomorphic to \mathbf{R}^n) and $U_T = (0,T] \times U$ for some fixed T > 0. We say that the partial differential operator $\frac{\partial}{\partial t} + L$ is parabolic if there exists a constant $\theta > 0$ such that $\sum_{i,j=1}^n a^{ij}(t,x)\xi_i\xi_j \ge \theta|\xi|^2$ for all $(t,x) \in U_T, \xi \in \mathbf{R}^n$. The letter L denotes for each time t a second order partial differential operator, having either the divergence form $Lu = -\sum_{i,j=1}^n (a^{ij}(t,x)u_{x_i})_{x_j} + \sum_{i=1}^n b^i(t,x)u_{x_i} + c(t,x)u$ or else the non-divergence form $Lu = -\sum_{i,j=1}^n a^{ij}(t,x)u_{x_ix_j} + \sum_{i=1}^n b^i(t,x)u_{x_i} + c(t,x)u$, for given coefficients a^{ij}, b^i, c (i, j = 1, 2, ..., n).

A family{ $G(t) : t \in \mathbf{R}_+$ } of bounded linear operators from X into E is a C_0 -semigroup (also called linear semigroup of class (C_0)) on X if (i) G(0) = the identity operator, and $G(t + s) = G(t)G(s) \forall t, s \ge 0$; (ii) $G(\cdot)$ is strongly continuous in $t \in \mathbf{R}_+$;

(iii) $||G(t)|| \leq Me^{\omega t}$ for some M > 0, real ω and $t \in \mathbf{R}_+$.

Example 4.2. Set $k(t,\tau)g(t,\tau,u) = G(t-\tau)u$, $\Phi(x) = x$, $\lambda(t) = G(t)x_0$, where $\{G(t)\}_{t\geq 0}$ is a C_0 -semigroup with an infinitesimal generator A. Then a solution of system (4.1)-(4.2) represents a mild solution of

$$x'(t) \in Ax(t) + F(t, x(t)), \qquad x(0) = x_0.$$
(4.3)

In particular, this problem includes control systems governed by parabolic partial differential equations as a special case. When A = 0, the relation (4.3) reduces to

$$x'(t) \in F(t, x(t)), \qquad x(0) = x_0.$$
 (5.4)

Denote

$$\Phi(u)(t) = \int_0^T k(t,\tau)g(t,\tau,u(\tau)) \, d\tau, \ t \in I.$$
(4.5)

Then the integral inclusion system (4.1)-(4.2) reduces to the form

$$x(t) = \lambda(t) + \Phi(u)(t) \qquad a.e. (I), \tag{S}$$

which may be written in more "compact" form as

$$u(t) \in F(t, V(\lambda + \Phi(u))(t))$$
 a.e. (I).

Now we recall the following:

Definition 4.3. A pair of functions (x, u) is called a solution pair of integral inclusion system (S), if $x(\cdot) \in C(I, E), u(\cdot) \in L^p(I, E)$ and satisfy relation (S).

For our further discussion, we denote by $S(\lambda)$ the solution set of (4.1) - (4.2).

For given $\alpha \in \mathbf{R}$ we denote by $L^p(I, E)$ the Banach space of all Bochner integrable functions $u(\cdot): I \to E$ endowed with the norm

$$||u(\cdot)||_p = \left(\int_0^T e^{-\alpha M_1 M_2 M_3 m(t)} ||u(t)||^p \, dt\right)^{\frac{1}{p}},$$

where $m(t) = \int_0^t L(s) \, ds$, $t \in I$. For our further discussion, we denote L = m(T).

Theorem 4.4. Let Hypothesis 4.1 be satisfied, let $\lambda(\cdot), \mu(\cdot) \in C(I, E)$ and let $v(\cdot) \in L^p(I, E)$ be such that

$$d(v(t), F(t, V(y)(t))) \le p(t) \qquad a.e. \qquad (I)$$

where $p(\cdot) \in L^p(I, \mathbf{R}_+)$ and $y(t) = \mu(t) + \Phi(v)(t), \forall t \in I$. Then for every $\alpha > 1$, there exists $x(\cdot) \in S(\lambda)$ such that for every $t \in I$

$$\begin{aligned} \|x(t) - y(t)\| &\leq \|\lambda - \mu\|_{C} + M_{1}M_{3}e^{\alpha M_{1}M_{2}M_{3}L} \Big[\frac{1}{\alpha^{\frac{1}{2p}}(\alpha^{\frac{1}{2p}} - 1)M_{1}^{\frac{1}{p}}M_{3}^{\frac{1}{p}}} \|\lambda - \mu\|_{C} \\ &+ \frac{\alpha^{\frac{1}{2p}}}{\alpha^{\frac{1}{2p}} - 1} \Big(\int_{0}^{T} e^{-\alpha M_{1}M_{2}M_{3}m(t)}p(t)dt\Big)^{\frac{1}{p}}\Big]^{p}. \end{aligned}$$

Proof. For $\lambda \in C(I, E)$ and $u \in L^p(I, E)$, define

$$x_{u,\lambda}(t) = \lambda(t) + \int_0^T k(t,s) g(t,s,u(s)) \, ds, \, t \in I.$$

Let us consider that $\lambda \in C(I, E), \sigma \in L^p(I, E)$ and define the set-valued maps

$$M_{\lambda,\sigma}(t) := F(t, V(x_{\sigma,\lambda})(t)), \ t \in I,$$
(4.6)

$$T_{\lambda}(\sigma) := \{\psi(\cdot) \in L^p(I, E) : \psi(t) \in M_{\lambda,\sigma}(t) \quad a.e. \ (I)\}.$$

$$(4.7)$$

Further, in view of condition (C3) of Hypothesis 4.1(H_3), $T_{\lambda}(\cdot)$ satisfies the condition: For any $\sigma \in L^p(I, E)$, $\sigma_1 \in T_{\lambda}(\sigma)$ and any given $\epsilon > 0$ there exists $\sigma_2 \in T_{\lambda}(\sigma_1)$ such that

$$\|\sigma_1 - \sigma_2\|_p \le H^+(T_\lambda(\sigma), T_\lambda(\sigma_1)) + \epsilon.$$
(4.8)

Now we claim that $T_{\lambda}(\sigma)$ is nonempty, bounded and closed for every $\sigma \in L^p(I, E)$.

It is well known that the set-valued map $M_{\lambda,\sigma}(\cdot)$ is measurable. For example the map $t \to M_{\lambda,\sigma}(t)$ can be approximated by step functions and so we can apply Theorem III. 40 in [1]. As the values of F are closed, with the measurable selection theorem we infer that $M_{\lambda,\sigma}(\cdot)$ is nonempty.

Further, we note that the set $T_{\lambda}(\sigma)$ is bounded and closed. Indeed, if $\psi_n \in T_{\lambda}(\cdot)$ and $\|\psi_n - \psi\|_p \to 0$, then there exists a subsequence ψ_{n_k} such that $\psi_{n_k}(t) \to \psi(t)$ for a.e. $t \in I$ and we find that $\psi \in T_{\lambda}(\sigma)$.

Let $\sigma_1, \sigma_2 \in L^p(I, E)$ be given. Let $\psi_1 \in T_{\lambda}(\sigma_1)$ and let $\delta > 0$. Consider the following set-valued map:

$$\mathcal{G}(t) := M_{\lambda,\sigma_2}(t) \cap \Big\{ z \in E : \|\psi_1(t) - z\|^p \le M_1 M_2 M_3 L(t) \int_0^T \|\sigma_1(s) - \sigma_2(s)\|^p \, ds + \delta \Big\}.$$

By (C2), it follows that

$$d^{p}(\psi_{1}(t), M_{\lambda,\sigma_{2}}(t)) \leq H^{+} \Big(F(t, V(x_{\sigma_{1},\lambda})(t)), F(t, V(x_{\sigma_{2},\lambda})(t)) \Big) + \epsilon$$

$$\leq L(t) \| V(x_{\sigma_{1},\lambda})(t)) - V(x_{\sigma_{2},\lambda})(t)) \| + \epsilon$$

$$\leq M_{2}L(t) \| x_{\sigma_{1},\lambda}(t) - x_{\sigma_{2},\lambda}(t) \| + \epsilon$$

$$\leq M_{2}M_{3}L(t) \int_{0}^{T} \| g(t,s,\sigma_{1}(s)) - g(t,s,\sigma_{2}(s)) \| ds + \epsilon$$

$$\leq M_{1}M_{2}M_{3}L(t) \int_{0}^{T} \| \sigma_{1}(s) - \sigma_{2}(s) \|^{p} ds + \epsilon.$$

Since ϵ is arbitrary, letting $\epsilon \to 0$, we deduce that $\mathcal{G}(\cdot)$ is nonempty bounded and has closed values. Further, according to Proposition III.4 in [1], $\mathcal{G}(\cdot)$ is measurable.

Let $\psi_2(\cdot)$ be a measurable selector of $\mathcal{G}(\cdot)$. It follows that $\psi_2 \in T_\lambda(\sigma_2)$ and

$$\begin{split} \|\psi_1 - \psi_2\|_p^p &= \int_0^T e^{-\alpha M_1 M_2 M_3 m(t)} \|\psi_1(t) - \psi_2(t)\|^p dt \\ &\leq \int_0^T e^{-\alpha M_1 M_2 M_3 m(t)} (M_1 M_2 M_3 L(t) \int_0^T \|\sigma_1(s) - \sigma_2(s)\|^p ds) dt \\ &+ \delta \int_0^T e^{-\alpha M_1 M_2 M_3 m(t)} dt \\ &\leq \frac{1}{\alpha} \|\sigma_1 - \sigma_2\|_p^p + \delta \int_0^T e^{-\alpha M_1 M_2 M_3 m(t)} dt. \end{split}$$

Since δ is arbitrary, so letting $\delta \to 0$ we deduce from the above inequality that

$$\|\psi_1 - \psi_2\|_p^p \le \frac{1}{\alpha} \|\sigma_1 - \sigma_2\|_p^p$$

i.e.,

$$\|\psi_1 - \psi_2\|_p \le \frac{1}{\alpha^{\frac{1}{p}}} \|\sigma_1 - \sigma_2\|_p$$

This yields

$$d(\psi_1, T_\lambda(\sigma_2)) \le \frac{1}{\alpha^{\frac{1}{p}}} \|\sigma_1 - \sigma_2\|_p$$

Thus, we have

$$\rho(T_{\lambda}(\sigma_1), T_{\lambda}(\sigma_2)) = \sup_{\psi_1 \in T_{\lambda}(\sigma_1)} d(\psi_1, T_{\lambda}(\sigma_2)) \le \frac{1}{\alpha^{\frac{1}{p}}} \|\sigma_1 - \sigma_2\|_p.$$
(4.9)

Now replacing $\sigma_1(\cdot)$ with $\sigma_2(\cdot)$ and arguing as above, we obtain

$$\rho(T_{\lambda}(\sigma_2), T_{\lambda}(\sigma_1)) \le \frac{1}{\alpha^{\frac{1}{p}}} \|\sigma_1 - \sigma_2\|_p.$$

$$(4.10)$$

Now adding (4.9) and (4.10) and dividing by 2, we obtain

$$H^{+}(T_{\lambda}(\sigma_{1}), T_{\lambda}(\sigma_{2})) \leq \frac{1}{\alpha^{\frac{1}{p}}} \|\sigma_{1} - \sigma_{2}\|_{p}$$

$$\leq \frac{1}{\alpha^{\frac{1}{p}}} \max\{\|\sigma_{1} - \sigma_{2}\|_{p}, d(\sigma_{1}, T_{\lambda}(\sigma_{1})), d(\sigma_{2}, T_{\lambda}(\sigma_{2})), [d(\sigma_{1}, T_{\lambda}(\sigma_{2})) + d(\sigma_{2}, T_{\lambda}(\sigma_{1}))]/2\}.$$

Hence we conclude that $T_{\lambda}(\cdot)$ is an H^+ -type multi-valued weak contractive mapping on $L^p(I, E)$. Next, we consider the following set-valued maps

$$F(t,x) := F(t,x) + p(t),$$

$$\tilde{M}_{\lambda,\sigma}(t) := \tilde{F}(t, V(x_{\sigma,\lambda})(t)), \qquad t \in I,$$

$$\tilde{T}_{\lambda}(\sigma) := \{\psi(\cdot) \in L^{p}(I,E) : \psi(t) \in \tilde{M}_{\lambda,\sigma}(t) \text{ a.e. } (I)\}$$

It is obvious that $\tilde{F}(\cdot, \cdot)$ satisfies Hypothesis 4.1. Let $\phi \in T_{\lambda}(\sigma), \delta > 0$ and define

$$\mathcal{G}_1(t) := \tilde{M}_{\lambda,\sigma}(t) \cap \Big\{ z \in X : \|\phi(t) - z\|^p \le M_2 L(t) \|\lambda - \mu\|_C^p + p(t) + \delta \Big\}.$$

Using the same argument as used for the set valued map $\mathcal{G}(\cdot)$, we deduce that $\mathcal{G}_1(\cdot)$ is measurable with nonempty closed values.

Next, we prove the following estimate:

$$H^{+}(T_{\lambda}(\sigma), \tilde{T}_{\mu}(\sigma)) \leq \frac{1}{\alpha^{\frac{1}{p}} M_{1}^{\frac{1}{p}} M_{3}^{\frac{1}{p}}} \|\lambda - \mu\|_{C} + \left(\int_{0}^{T} e^{-\alpha M_{1} M_{2} M_{3} m(t)} p(t) dt\right)^{\frac{1}{p}}.$$
 (4.11)

Let $\psi(\cdot) \in \tilde{T}_{\mu}(\sigma)$. Then

$$\begin{split} \|\phi - \psi\|_{p}^{p} &= \int_{0}^{T} e^{-\alpha M_{1}M_{2}M_{3}m(t)} \|\phi(t) - \psi(t)\|^{p} dt \\ &\leq \int_{0}^{T} e^{-\alpha M_{1}M_{2}M_{3}m(t)} [M_{2}L(t)\|\lambda - \mu\|_{C}^{p} + p(t) + \delta] dt \\ &\leq \|\lambda - \mu\|_{C}^{p} \int_{0}^{T} e^{-\alpha M_{1}M_{2}M_{3}m(t)} M_{2}L(t) dt \\ &+ \int_{0}^{T} e^{-\alpha M_{1}M_{2}M_{3}m(t)} p(t) dt + \delta \int_{0}^{T} e^{-\alpha M_{1}M_{2}M_{3}m(t)} dt \\ &\leq \frac{1}{\alpha M_{1}M_{3}} \|\lambda - \mu\|_{C}^{p} + \int_{0}^{T} e^{-\alpha M_{1}M_{2}M_{3}m(t)} p(t) dt \\ &+ \delta \int_{0}^{T} e^{-\alpha M_{1}M_{2}M_{3}m(t)} dt. \end{split}$$

Since δ is arbitrary, so letting $\delta \to 0$ we deduce from the above inequality that

$$\|\phi - \psi\|_p^p \le \frac{1}{\alpha M_1 M_3} \|\lambda - \mu\|_C^p + \int_0^T e^{-\alpha M_1 M_2 M_3 m(t)} p(t) dt.$$

Thus, by taking $\frac{1}{p}$ th power on both sides of the above inequality breaking the right hand side, one obtains (4.11).

Now applying Proposition 3.8 we obtain

$$H^{+}(Fix(T_{\lambda}), Fix(\tilde{T}_{\mu})) \leq \frac{1}{\alpha^{\frac{1}{2p}}(\alpha^{\frac{1}{2p}} - 1)M_{1}^{\frac{1}{p}}M_{3}^{\frac{1}{p}}} \|\lambda - \mu\|_{C} + \frac{\alpha^{\frac{1}{2p}}}{\alpha^{\frac{1}{2p}} - 1} \Big(\int_{0}^{T} e^{-\alpha M_{1}M_{2}M_{3}m(t)}p(t)dt\Big)^{\frac{1}{p}}.$$

Since $v(\cdot) \in Fix(\tilde{T}_{\mu})$, it follows that there exists $u(\cdot) \in Fix(T_{\lambda})$ such that

$$\|v - u\|_{p} \leq \frac{1}{\alpha^{\frac{1}{2p}} (\alpha^{\frac{1}{2p}} - 1)M_{1}^{\frac{1}{p}} M_{3}^{\frac{1}{p}}} \|\lambda - \mu\|_{C} + \frac{\alpha^{\frac{1}{2p}}}{\alpha^{\frac{1}{2p}} - 1} \Big(\int_{0}^{T} e^{-\alpha M_{1} M_{2} M_{3} m(t)} p(t) dt\Big)^{\frac{1}{p}}.$$
 (4.12)

We define

$$x(t) = \lambda(t) + \int_0^T k(t,s) g(t,s,u(s)) \, ds.$$

Then one has the following inequality:

$$\|x(t) - y(t)\| \le \|\lambda(t) - \mu(t)\| + M_1 M_3 \int_0^T \|u(s) - v(s)\|^p ds$$

$$\le \|\lambda - \mu\|_C + M_1 M_3 e^{\alpha M_1 M_2 M_3 L} \|u - v\|_p^p.$$

Combining the last inequality with (4.12) we obtain

$$\begin{aligned} \|x(t) - y(t)\| &\leq \|\lambda - \mu\|_{C} + M_{1}M_{3}e^{\alpha M_{1}M_{2}M_{3}L} \Big[\frac{1}{\alpha^{\frac{1}{2p}}(\alpha^{\frac{1}{2p}} - 1)M_{1}^{\frac{1}{p}}M_{3}^{\frac{1}{p}}} \|\lambda - \mu\|_{C} \\ &+ \frac{\alpha^{\frac{1}{2p}}}{\alpha^{\frac{1}{2p}} - 1} \Big(\int_{0}^{T} e^{-\alpha M_{1}M_{2}M_{3}m(t)}p(t)dt\Big)^{\frac{1}{p}}\Big]^{p}. \end{aligned}$$

This completes the proof.

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