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# DISCONTINUOUS WAVE EQUATIONS AND A TOPOLOGICAL DEGREE FOR SOME CLASSES OF MULTI-VALUED MAPPINGS 

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Abstract. The Leray-Schauder degree is extended to certain multi-valued mappings on separable Hilbert spaces with applications to the existence of weak periodic solutions of discontinuous semilinear wave equations with fixed ends.

Keywords: discontinuous wave equations, topological degree, multi-valued mappings
MSC 2000: 35L05, 47H17, 58C06

## 1. INTRODUCTION

In this paper we study the existence of weak $2 \pi$-periodic solutions to the discontinuous semilinear wave equation

$$
\begin{gather*}
u_{t t}-u_{x x}+g(u)+f(x, t, u)=h(x, t),  \tag{Prob}\\
u(0, \cdot)=u(\pi, \cdot)=0
\end{gather*}
$$

where $f: \Omega \times \mathbb{R} \rightarrow \mathbb{R}, \Omega=(0, \pi) \times(0,2 \pi)$ is Carathéodory continuous and nondecreasing in $u, g: \mathbb{R} \rightarrow \mathbb{R}$ is bounded nondecreasing and $h \in L^{2}(\Omega)$. Moreover, we suppose

$$
|f(x, t, u)| \leqslant c_{0}|u|+h_{0}(x, t) \quad \forall u \in \mathbb{R}, \forall(x, t) \in \Omega
$$

for a constant $c_{0}>0$ and $h_{0} \in L^{2}(\Omega)$. (Prob) with $g=0$, i.e. for the continuous case, was studied in [2], where a construction of a topological degree is introduced for a class of monotone single-valued mappings. The main purpose of this paper is to extend that method to monotone multi-valued mappings. Like in [2], we use a Galerkin projection method, which is however not at all the standard Galerkin approximation method of [2]. We construct continuous one-parametric generalized

Galerkin projections which we use for the derivation of a one-parametric family of multi-valued mappings possessing the Leray-Schauder degree [7]. The basic Lemma 6 below is the stabilization of this degree for large parameters. In this way, we can define a topological degree for our multi-valued mappings. The rest of the paper is the extension of the main results of [2] to the multi-valued case. We end the paper with a discontinuous equation of a vibrating string presenting a problem at resonance with an infinite-dimensional kernel. This result is motivated by [4]. Other topological degrees for single and multi-valued mappings have been introduced in [1], [3], [5], [6] and [9].

## 2. Standard classes of mappings

Let $H$ be a real separable Hilbert space with an inner product $(\cdot, \cdot)$. In what follows we use the following notation:

- $\left\{a_{n}\right\}$ means $\left\{a_{n}\right\}_{n=1}^{\infty}$;
- $\lim a_{n}$ means $\lim _{n \rightarrow \infty} a_{n}$ (similarly limsup and $\lim \inf$ );
- | $\cdot$ always denotes the norm in $H$;
- $\|\cdot\|$ denotes the norm in the space of continuous linear mappings $\mathscr{L}(H)$.

A multi-valued mapping $F: H \rightarrow 2^{H}$ is

- monotone (denote $F \in(m M O N)$ ), if

$$
\left(f_{u}^{*}-f_{v}^{*}, u-v\right) \geqslant 0
$$

for all $u, v \in H$ and all selections $f_{u}^{*} \in F(u), f_{v}^{*} \in F(v)$;

- quasimonotone $(F \in(m Q M))$, if for any sequence $\left\{u_{n}\right\}$ in $H$ with $u_{n} \rightharpoonup u$ and for all selections $f_{n}^{*} \in F\left(u_{n}\right)$ we have

$$
\liminf \left(f_{n}^{*}, u_{n}^{*}-u\right) \geqslant 0
$$

- pseudomonotone $\left(F \in(m P M)\right.$ ), if for any sequence $\left\{u_{n}\right\}$ in $H$ with $u_{n} \rightharpoonup u$, the existence of selections $f_{n}^{*} \in F\left(u_{n}\right)$ and a point $f^{*} \in H$ with $f_{n}^{*} \rightharpoonup f^{*}$ and $\lim \sup \left(f_{n}^{*}, u_{n}-u\right) \leqslant 0$ implies that $f^{*} \in F(u)$ and $\left(f_{n}^{*}, u_{n}\right) \rightarrow\left(f^{*}, u\right)$;
- of class $\left(m S_{+}\right)\left(F \in\left(m S_{+}\right)\right)$, if for any sequence $\left\{u_{n}\right\}$ in $H$ with $u_{n} \rightharpoonup u$, the existence of selections $f_{n}^{*} \in F\left(u_{n}\right)$ with $\lim \sup \left(f_{n}^{*}, u_{n}-u\right) \leqslant 0$ implies that $u_{n} \rightarrow u$;
- compact $(F \in(m C O M P))$, if for any bounded sequence $\left\{u_{n}\right\}$ in $H$ and for any $f_{n}^{*} \in F\left(u_{n}\right)$ the sequence $\left\{f_{n}^{*}\right\}$ has a convergent subsequence;
- of Leray-Schauder type $(F \in(m L S))$, if $F=I+C$, where $I$ is the identity on $H$, for some $C \in(m C O M P)$;
- bounded, if for any bounded set $B \subset H$ the set $\bigcup_{u \in B} F(u)$ is bounded;
- convex-valued, if $F(u)$ is a non-empty convex set in $H$ for any $u \in H$;
- upper semicontinuous ( $F$ is usc), if $F^{-1}(A)$ is closed in $H$ whenever $A \subset H$ is closed;
- weak upper semicontinuous ( $F$ is w-usc), if for any sequence $\left\{u_{n}\right\} \in H, u_{n} \rightarrow$ $u \in H$, the existence of selections $f_{n}^{*} \in F\left(u_{n}\right)$ with $f_{n}^{*} \rightharpoonup f^{*} \in H$ implies $f^{*} \in F(u)$.
In what follows, we assume that all mappings used are bounded, w-usc and convexvalued. When a mapping is defined only on a subset of $H$, the above definitions can be modified in an obvious way.

Proposition 1. For the classes defined above the following inclusions hold:

$$
\begin{gathered}
(m L S) \subset\left(m S_{+}\right) \subset(m P M) \subset(m Q M) \\
(m M O N) \subset(m P M), \quad(m C O M P) \subset(m Q M)
\end{gathered}
$$

Proposition 2. If $F \in(m Q M)$ and $G \in\left(m S_{+}\right)$then $F+G \in\left(m S_{+}\right)$.
Proofs of Propositions 1 and 2 are similar to those for single-valued mappings [2].

## 3. Classes of admissible mappings

Let $G$ be a bounded open subset in $H, M$ a closed subspace of $H$ and let $Q$ and $P$ stand for the orthogonal projections to $M$ and $M^{\perp}$, respectively. The family

$$
\mathscr{F}_{G}=\left\{F: \bar{G} \rightarrow 2^{H} \mid F=Q g+P f \text { for some } g \in(m L S) \text { and } f \in\left(m S_{+}\right)\right\}
$$

is called the class of admissible mappings. Other very important and more general classes are
$\mathscr{F}_{G}(m Q M)=\left\{F: \bar{G} \rightarrow 2^{H} \mid F=Q g+P f\right.$ for some $g \in(m L S)$ and $\left.f \in(m Q M)\right\}$
and $\mathscr{F}_{G}(m P M)$ defined accordingly. Obviously $\mathscr{F}_{G}=\mathscr{F}_{G}\left(m S_{+}\right) \subset \mathscr{F}_{G}(m P M) \subset$ $\mathscr{F}_{G}(m Q M)$.

Let $L: H \supset D(L) \rightarrow H$ be a closed densely defined linear operator with $\operatorname{Im} L=$ $(\operatorname{Ker} L)^{\perp}$. Let $L_{0}=L / \operatorname{Im} L$ and assume that the right inverse $L_{0}^{-1}: \operatorname{Im} L \rightarrow \operatorname{Im} L$ is compact. We choose $M=\operatorname{Im} L$ and $M^{\perp}=\operatorname{Ker} L$. Let $N: H \rightarrow 2^{H}$. Then similarly to [2], we consider the mapping

$$
F=Q\left(I+L_{0}^{-1} Q N\right)+P N .
$$

Clearly, $F \in \mathscr{F}_{G}$ for $N \in\left(m S_{+}\right)$.

Lemma 3. Let $F$ and $N$ be defined as above. Then for $h \in H$

$$
\begin{equation*}
h \in L u+N(u) \text { with } u \in D(L) \cap \bar{G} \tag{SInc}
\end{equation*}
$$

if and only if

$$
y \in F(u) \text { with } u \in \bar{G} \text { for } y=\left(L_{0}^{-1} Q+P\right) h
$$

The proof is the same as for single-valued mappings [2]. The inclusion $h \in L u-$ $N(u)$ can be treated analogously.

## 4. Construction of the degree for $\mathscr{F}_{G}\left(m S_{+}\right)$

We construct a topological degree function for $\mathscr{F}_{G}$. First, we define a class of admissible homotopies. A mapping: $(t, u) \rightarrow f_{t}(u)$ from $[0,1] \times \bar{G}$ to $2^{H}$ is a (multi-) homotopy of the class $\left(m S_{+}\right)$, if for any sequences $\left\{u_{n}\right\}$ in $\bar{G},\left\{t_{n}\right\}$ in $[0,1], f_{n}^{*} \in$ $f_{t_{n}}\left(u_{n}\right)$ with $u_{n} \rightharpoonup u, t_{n} \rightarrow t$ and $\lim \sup \left(f_{n}^{*}, u_{n}-u\right) \leqslant 0$ we have $u_{n} \rightarrow u$. We also assume that the mapping $(t, u) \rightarrow f_{t}(u)$ satisfies all the necessary conditions.

Similarly, a mapping: $(t, u) \rightarrow g_{t}(u)=\left(I+C_{t}\right)(u)$ from $[0,1] \times \bar{G}$ to $2^{H}$ is called a (multi-)homotopy of the Leray-Schauder type, if the mapping $(t, u) \rightarrow C_{t}(u)$ is compact. We denote

$$
\mathscr{H}_{G}=\left\{F_{t} \mid F_{t}=Q g_{t}+P f_{t}\right\}
$$

where $g_{t}$ and $f_{t}$ are homotopies of the Leray-Schauder type and of the class $\left(m S_{+}\right)$, respectively. The set $\mathscr{H}_{G}$ is called the class of admissible homotopies. Obviously $F_{t}=(1-t) F_{1}+t F_{2} \in \mathscr{H}_{G}, 0 \leqslant t \leqslant 1$ for any $F_{1}, F_{2} \in \mathscr{F}_{G}$.

We use Galerkin approximations [8] with respect to the subspace $M^{\perp}$. The space $H\left(M^{\perp} \subset H\right)$ is separable so there exists a sequence $\left\{N_{n}\right\}$ of finite dimensional subspaces of $M^{\perp}$ with $N_{n} \subset N_{n+1}$ for all $n$, and $\cup_{n=1}^{\infty} N_{n}$ is dense in $M^{\perp}$. We denote by $P_{n}$ the orthogonal projection from $H$ to $N_{n}$. We extend this to generalized Galerkin approximations defined by

$$
P_{\lambda}=(\lambda-n) P_{n+1}+(n+1-\lambda) P_{n} \quad \text { for any } \lambda \in[n, n+1] .
$$

We have the following obvious result.

Proposition 4. The generalized Galerkin approximations satisfy
i) $\left(P_{\lambda} u, v\right)=\left(u, P_{\lambda} v\right)$ for every $\lambda \geqslant 1, u, v \in H$;
ii) $\left\|P_{\lambda}\right\| \leqslant 1$ for all $\lambda \geqslant 1$;
iii) $P_{\lambda} v \rightarrow P v$ for every $v \in H$ as $\lambda \rightarrow \infty$;
iv) $P_{n} P_{\lambda}=P_{n}$ for every $\lambda \geqslant n \in \mathbb{N}$;
v) $\left(z, P_{\lambda} z\right) \geqslant 0$ for every $z \in H$ and $\lambda \geqslant 1$.

For each $F=Q(I+C)+P f \in \mathscr{F}_{G}$, we define the approximations $\left\{F_{\lambda} \mid \lambda \geqslant 1\right\}$ by

$$
F_{\lambda}=I+Q C+\lambda P_{\lambda} f
$$

We note that $Q C+\lambda P_{\lambda} f$ is compact, convex-valued and usc for each $\lambda \geqslant 1$.
Similarly, for each admissible homotopy $F_{t}=Q\left(I+C_{t}\right)+P f_{t}, 0 \leqslant t \leqslant 1$, we have

$$
\left(F_{t}\right)_{\lambda}=I+Q C_{t}+\lambda P_{\lambda} f_{t}
$$

which is obviously a homotopy of the Leray-Schauder type for any $\lambda \geqslant 1$. Finally, if $y \in H$ is given, we denote

$$
y_{\lambda}=Q y+\lambda P_{\lambda} y \quad \text { for each } \lambda \geqslant 1 .
$$

It is clear that $(F-y)_{\lambda}=F_{\lambda}-y_{\lambda}$.

Proposition 5. Let $\left\{u_{k}\right\}, u_{k} \in H$ and let $\left\{P_{\lambda}\right\}$ be the projections defined as above. Then
a) if $u_{k} \rightharpoonup u$ and $\lambda_{k} \rightarrow \infty$ then $P_{\lambda_{k}} u_{k} \rightharpoonup P u$;
b) if $u_{k} \rightarrow u$ and $\lambda_{k} \rightarrow \infty$ then $P_{\lambda_{k}} u_{k} \rightarrow P u$.

Now we can formulate the basic lemma.

Lemma 6. Let $F_{t}=Q\left(I+C_{t}\right)+P f_{t}$ be an admissible homotopy, $y_{t}(0 \leqslant t \leqslant 1)$ a continuous curve in $H$ and let $A$ be a closed subset in $\bar{G}$. If $y_{t} \notin F_{t}(A)$ for all $t \in[0,1]$, then there exists $n_{0} \in \mathbb{N}$ such that

$$
\left(y_{t}\right)_{\lambda} \notin\left(F_{t}\right)_{\lambda}(A) \quad \text { for all } t \in[0,1] \text { and } \lambda \geqslant n_{0} .
$$

Proof. Since also $F_{t}-y_{t}$ defines an admissible homotopy and $\left(F_{t}-y_{t}\right)_{\lambda}=$ $\left(F_{t}\right)_{\lambda}-\left(y_{t}\right)_{\lambda}$, we may assume, without loss of generality, that $y_{t} \equiv 0$.

If the assertion were false, there would exist sequences $\left\{u_{k}\right\}$ in $A$ and $\left\{\lambda_{k}\right\}$ in $[1, \infty), \lambda_{k} \rightarrow \infty$, and $\left\{t_{k}\right\}$ in $[0,1]$ such that $0 \in\left(F_{t_{k}}\right)_{\lambda_{k}}\left(u_{k}\right)$. This is equivalent to the existence of selections $g_{k}^{*} \in f_{t_{k}}\left(u_{k}\right)$ and $c_{k}^{*} \in C_{t_{k}}\left(u_{k}\right)$ for which

$$
u_{k}+Q c_{k}^{*}+\lambda_{k} P_{\lambda_{k}} g_{k}^{*}=0
$$

Writing this equation in both subspaces $M$ and $M^{\perp}$ we get

$$
\begin{align*}
Q u_{k}+Q c_{k}^{*} & =0,  \tag{1}\\
P u_{k}+\lambda_{k} P_{\lambda_{k}} g_{k}^{*} & =0 . \tag{2}
\end{align*}
$$

The sequence $u_{k}$ is bounded therefore we can (taking a subsequence, if necessary) assume that $u_{k} \rightharpoonup u$ for some $u \in H$. Similarly we can assume $t_{k} \rightarrow t, t \in[0,1]$. We have $c_{k}^{*} \in C_{t_{k}}\left(u_{k}\right)$. Since $C$ is a compact mapping we can assume $c_{k}^{*} \rightarrow z^{*}, z^{*} \in H$ and, since the sequence $g_{k}^{*}$ is bounded, $g_{k}^{*} \rightharpoonup g^{*}, g^{*} \in H$.

Hence we have $Q c_{k}^{*} \rightarrow Q z^{*}$ and by (1), $Q u_{k} \rightarrow Q u$. By (2),

$$
\frac{1}{\lambda_{k}} P u_{k}+P_{\lambda_{k}} g_{k}^{*}=0 .
$$

The set $\left\{u_{k} \mid k=1,2, \ldots\right\}$ is bounded, thus $\frac{1}{\lambda_{k}} P u_{k} \rightarrow 0$ for $k \rightarrow \infty$. This leads to $P_{\lambda_{k}} g_{k}^{*} \rightarrow 0$. On the other hand, for $g_{k}^{*} \rightharpoonup g^{*}$ we conclude $P_{\lambda_{k}} g_{k}^{*} \rightharpoonup P g^{*}$, which yields $P g^{*}=0$. Hence we have $P g_{k}^{*} \rightharpoonup P g^{*}=0$ followed by $\lim \left(g_{k}^{*}, P u\right)=0, \forall u \in H$.

We continue with calculating of $\lim \sup \left(g_{k}^{*}, u_{k}-u\right)$ :

$$
\begin{equation*}
\limsup \left(g_{k}^{*}, u_{k}-u\right)=\lim \sup \left(g_{k}^{*}, P u_{k}-P u\right)=\limsup \left(g_{k}^{*}, P u_{k}\right) . \tag{3}
\end{equation*}
$$

From (2) we obtain $P u_{k}=-\lambda_{k} P_{\lambda_{k}} g_{k}^{*}$. Inserting it in to the last term of (3) we get

$$
\limsup \left(g_{k}^{*}, u_{k}-u\right)=-\liminf \lambda_{k}\left(g_{k}^{*}, P_{\lambda_{k}} g_{k}^{*}\right)
$$

By v) of Proposition $4, \lambda_{k}\left(g_{k}^{*}, P_{\lambda_{k}} g_{k}^{*}\right) \geqslant 0$ and it immediately follows that

$$
\limsup \left(g_{k}^{*}, u_{k}-u\right) \leqslant 0
$$

Hence, by the definition of the homotopy $f_{t}$ of the class $\left(m S_{+}\right)$we have $u_{k} \rightarrow u$, $u \in A$. We have $t_{n} \rightarrow t, g_{k}^{*} \rightharpoonup g^{*}$ and $u_{k} \rightarrow u$. Since $(t, u) \rightarrow f_{t}(u)$ is w-usc, we get $g^{*} \in f_{t}(u)$. Similarly, we get $z^{*} \in C_{t}(u)$. From $Q u_{k} \rightarrow Q u, Q c_{k}^{*} \rightarrow Q z^{*}$ using (1), we obtain $Q u+Q z^{*}=0$. Since $P g^{*}=0$ we have

$$
0=Q u+Q z^{*}+P g^{*} \in Q u+Q C_{t}(u)+P f_{t}(u)=F_{t}(u),
$$

a contradiction. The proof is complete.

If we choose a constant homotopy $F_{t}=F=Q g+P f, y_{t}=y$, and $A=\partial G$, we obtain the stabilization of a degree. If $y \notin F(\partial G)$, then there exists $\lambda_{0} \geqslant 1$ such that $y_{\lambda} \notin F_{\lambda}(\partial G)$ for all $\lambda \geqslant \lambda_{0}$.

Lemma 7. Let $F \in \mathscr{F}_{G}$ and $y \notin F(\partial G)$. Then there exists $\lambda_{1} \in[1, \infty)$ such that

$$
d\left(F_{\lambda}, G, y_{\lambda}\right)=\text { constant } \quad \text { for all } \lambda \geqslant \lambda_{1} .
$$

Due to Lemma 7 we can define a degree function for the class $\mathscr{F}_{G}$. We put

$$
\begin{equation*}
d(F, G, y)=\lim _{\lambda \rightarrow \infty} d\left(F_{\lambda}, G, y_{\lambda}\right) \tag{4}
\end{equation*}
$$

for any given $F \in \mathscr{F}_{G}$ and $y \in H$ with $y \notin F(\partial G)$. In the next lemma we show that the degree function defined by (4) has all the usual properties. The proof is the same as in [2] with the use of the degree theory for multi-valued mappings [7].

Lemma 8. Let $G$ be an open bounded subset of $H$ and $F \in \mathscr{F}_{G}\left(m S_{+}\right)$. Then
i) $d(F, G, y) \neq 0$ implies that there exists $y \in F(\bar{G})$.
ii) $d(F, G, y)=d\left(F, G_{1}, y\right)+d\left(F, G_{2}, y\right)$ (thus $F \in \mathscr{F}_{G_{1}}$ and $F \in \mathscr{F}_{G_{2}}$ ), whenever $G_{1}$ and $G_{2}$ are disjoint open subsets of $G$ such that $y \notin F\left(\bar{G} \backslash\left(G_{1} \cup G_{2}\right)\right)$.
iii) $d\left(F_{t}, G, y_{t}\right)$ is independent of $t \in J=[0,1]$ if $F_{t} \in \mathscr{H}_{G}, y_{t}$ is a continuous curve in $H$ and $y_{t} \notin F_{t}(\partial G)$ for all $t \in J$.
iv) $d(I, G, w)=1$ if and only if $w \in G$ (normalization).

We can also simply prove the Borsuk Theorem for the class $\mathscr{F}_{G}\left(m S_{+}\right)$.

Proposition 9. Let $G$ be an open bounded symmetric subset of $H$ containing the origin and let $F$ be a multi-mapping $\bar{G} \rightarrow 2^{H}$ satisfying $F(-u)=-F(u)$ for all $u \in \partial G$. Then if $F \in \mathscr{F}_{G}\left(m S_{+}\right)$, the inclusion $0 \in F(u)$ admits a solution $u$ in $\bar{G}$ and $d(F, G, 0)$ is odd whenever defined.

We have already proved all the desired properties of the degree function for the class $\mathscr{F}_{G}$. We formulate our result in the next theorem.

Theorem 10. Let $H$ be a real separable Hilbert space, $G$ a bounded open subset of $H, \mathscr{F}_{G}$ the class of admissible mappings and $\mathscr{H}_{G}$ the class of admissible homotopies defined above. Then there exists a classical topological degree function $d$ on $\mathscr{F}_{G}$ satisfying the properties i), ii), iii) and iv) from Lemma 8 with respect to $\mathscr{H}_{G}$ and the normalizing mapping $I$.

The most common method of studying the existence of solutions of differential equations (inclusions) applying a degree theory is to use a homotopy between a mapping and the normalizing mapping-the identity. However, it is not always possible to use such a homotopy. Sometimes we can substitute the identity with an other mapping with a nonzero degree. A mapping $R \in \mathscr{F}_{G}\left(m S_{+}\right)$with $d(R, G, y) \neq 0$ for all $y \in R(G)$ is called a reference mapping.

The next basic theorem immediately follows from the properties of the degree function $d$ (see [2, Theorem 2]).

Theorem 11. Let $G$ be an open bounded set, $G \subset H, R \in \mathscr{F}_{G}$ a reference mapping and $F \in \mathscr{F}_{G}$. If for a given $y \in H$ there exists $w \in R(G)$ such that

$$
(1-t) w+t y \notin(1-t) R(u)+t F(u) \quad \text { for all } u \in \partial G \text { and } t \in[0,1]
$$

then $d(F, G, y) \neq 0$ and hence the inclusion $y \in F(u)$ admits a solution $u$ in $G$.
We already have constructed a degree function for the class $\mathscr{F}_{G}=\mathscr{F}_{G}\left(m S_{+}\right)$. By using this we can also build a similar degree function for the larger class $\mathscr{F}_{G}(m Q M)$. By Proposition 2 we can use an $\left(m S_{+}\right)$approximation $f_{\varepsilon}=f+\varepsilon I$ for any given $f \in(m Q M)$ and $\varepsilon>0$. In fact, we generalize Theorem 11 in the following form (see [2, Theorem 3]).

Theorem 12. Let $G$ be an open bounded subset of $H, R \in \mathscr{F}_{G}\left(m S_{+}\right)$a reference mapping and $F \in \mathscr{F}_{G}(m Q M)$. If for a given $y \in H$ there exists $w \in R(G)$ such that the condition

$$
\begin{equation*}
(1-t) w+t y \notin(1-t) R(u)+t F(u) \quad \text { for all } u \in \partial G \text { and } t \in[0,1] \tag{5}
\end{equation*}
$$

holds, then the inclusion $y \in F(u)$ is almost solvable, i.e. $y \in \overline{F(\bar{G})}$.
Proof. Without loss of generality, we can assume that $y=0$ and $w=0 \in R(G)$. If $0 \in \overline{F(\partial G)} \subset \overline{F(\bar{G})}$, the assertion is true. Thus let $0 \notin \overline{F(\partial G)}$. Since $F \in$ $\mathscr{F}_{G}(m Q M), F$ has a representation $F=Q(I+C)+P f$ for some $C \in(m C O M P)$ and $f \in(m Q M)$. Similar argument leads to $R=Q\left(I+C_{0}\right)+P f_{0}$ for some $C_{0} \in$ $(m C O M P)$ and $f_{0} \in\left(m S_{+}\right)$. Now we introduce an $\mathscr{F}_{G}\left(m S_{+}\right)$approximation of $F$, for each $\varepsilon>0$ we denote

$$
F_{\varepsilon}=F+\varepsilon P R=Q(I+C)+P\left(f+\varepsilon f_{0}\right) .
$$

By Proposition $2, f+\varepsilon f_{0} \in\left(m S_{+}\right)$, therefore $F_{\varepsilon} \in \mathscr{F}_{G}\left(m S_{+}\right)$for all $\varepsilon>0$. For applying Theorem 11 to the mapping $F_{\varepsilon}$ we show there exists $\varepsilon_{0}$ such that
(6) $\quad 0 \notin(1-t) R(u)+t F_{\varepsilon}(u) \quad$ for all $u \in \partial G$ and $t \in[0,1]$ and $0<\varepsilon<\varepsilon_{0}$.

If (6) were not true, we could find sequences $\left\{\varepsilon_{n}\right\},\left\{t_{n}\right\} \subset[0,1]$ and $\left\{u_{n}\right\} \subset \partial G$ for which

$$
\begin{equation*}
0 \in\left(1-t_{n}\right) R\left(u_{n}\right)+t_{n} F_{\varepsilon_{n}}\left(u_{n}\right) \text { for all } n \in \mathbb{N} \tag{7}
\end{equation*}
$$

Taking subsequences, if necessary, we can assume that $\varepsilon_{n} \rightarrow 0^{+}, t_{n} \rightarrow t$ and $u_{n} \rightharpoonup u$. Clearly $t \neq 1$, since otherwise $0 \in \overline{F(\partial G)}$ which contradicts the assumption. Thus $t \in[0,1)$. Writing (6) in the components and taking selections $c_{n}^{*} \in C\left(u_{n}\right), d_{n}^{*} \in$ $C_{0}\left(u_{n}\right), f_{n}^{*} \in f\left(u_{n}\right)$ and $g_{n, 1}^{*}, g_{n, 2}^{*} \in f_{0}\left(u_{n}\right)$ we have

$$
\begin{array}{r}
Q u_{n}+\left(1-t_{n}\right) Q d_{n}^{*}+t_{n} Q c_{n}^{*}=0 \\
\left(1-t_{n}\right) P g_{n, 1}^{*}+t_{n} \varepsilon_{n} P g_{n, 2}^{*}+t_{n} P f_{n}^{*}=0 \tag{9}
\end{array}
$$

with $r_{n}^{*}=Q\left(u_{n}+d_{n}^{*}\right)+P g_{n, 1}^{*} \in R\left(u_{n}\right)$ and $h_{n}^{*}=Q\left(u_{n}+c_{n}^{*}\right)+P\left(f_{n}^{*}+\varepsilon_{n} g_{n, 2}^{*}\right) \in$ $F_{\varepsilon_{n}}\left(u_{n}\right)$. Without loss of generality, we can assume that $c_{n}^{*} \rightarrow c^{*}$ and $d_{n}^{*} \rightarrow d^{*}$ for some $c^{*}, d^{*} \in H$, which yields $r_{n}^{*} \rightharpoonup r^{*}$ and $h_{n}^{*} \rightharpoonup h^{*}$. Therefore from (8) we have $Q u_{n} \rightarrow Q u\left(u_{n} \rightharpoonup u\right)$. By (9)

$$
P g_{n, 1}^{*}=-\frac{t_{n}}{1-t_{n}} P f_{n}^{*}-\frac{t_{n} \varepsilon_{n}}{1-t_{n}} P g_{n, 2}^{*} .
$$

By using this and $Q u_{n} \rightarrow Q u$ we can calculate

$$
\begin{gathered}
\limsup \left(g_{n, 1}^{*}, u_{n}-u\right)=\lim \sup \left(P g_{n, 1}^{*}, u_{n}-u\right) \\
=\limsup \left\{-\frac{t_{n}}{1-t_{n}}\left(P f_{n}^{*}, u_{n}-u\right)-\frac{t_{n} \varepsilon_{n}}{1-t_{n}}\left(P g_{n, 2}^{*}, u_{n}-u\right)\right\} \\
=-\frac{t}{1-t} \liminf \left(f_{n}^{*}, u_{n}-u\right)
\end{gathered}
$$

Since $f \in(m Q M)$, we conclude that

$$
\limsup \left(g_{n, 1}^{*}, u_{n}-u\right)=-\frac{t}{1-t} \lim \inf \left(f_{n}^{*}, u_{n}-u\right) \leqslant 0
$$

By using the $\left(m S_{+}\right)$property of $f_{0}$ we obtain $u_{n} \rightarrow u$. Clearly $u \in \partial G$. From the wusc property of the mapping $R$ we get $r_{n}^{*} \rightharpoonup r^{*} \in R(u)$ and similarly $h_{n}^{*} \rightharpoonup h^{*} \in F(u)$. Consequently,

$$
0=\left(1-t_{n}\right) r_{n}^{*}+t_{n} h_{n}^{*} \rightharpoonup(1-t) r^{*}+t h^{*}
$$

giving $0=(1-t) r^{*}+t h^{*}$ and hence

$$
0 \in(1-t) R(u)+t F(u) \quad \text { for some } u \in \partial G
$$

That gives a contradiction to (5) (y=w=0). Hence we have proved (6). Thus applying Theorem 11 we find $u_{\varepsilon} \in G$ such that

$$
\begin{equation*}
0 \in F_{\varepsilon}\left(u_{\varepsilon}\right) \quad \text { for every } \varepsilon \in\left(0, \varepsilon_{0}\right) \tag{10}
\end{equation*}
$$

Choosing an arbitrary sequence $\left\{\varepsilon_{i}\right\}$ with $\varepsilon_{i} \rightarrow 0^{+}$along with solutions $u_{\varepsilon_{i}}$ from (10) we get

$$
F\left(u_{\varepsilon_{i}}\right)=F_{\varepsilon_{i}}\left(u_{\varepsilon_{i}}\right)-\varepsilon_{i} P R\left(u_{\varepsilon_{i}}\right)
$$

where the set $\cup_{i} P R\left(u_{\varepsilon_{i}}\right)$ is bounded. Thus $0 \in \overline{F(\bar{G})}$.
Because the condition of quasimonotonicity is in applications the easiest to verify (e.g. it is implied by the monotonicity), we need a result for $\mathscr{F}_{G}(m P M)$. By using similar argument as in Theorem 12, we obtain the following result (see [2, Theorem 4]).

Theorem 13. Let the assumptions of Theorem 12 be satisfied and let, in addition, $G$ be convex and $F \in \mathscr{F}_{G}(m P M)$. If for a given $y \in H$ there exists $w \in R(G)$ such that the condition (5) holds, then the inclusion $y \in F(u)$ does admit a solution $u$ in $G$.

Now we can generalize the Borsuk Theorem to the class $\mathscr{F}_{G}(m Q M)$ (see [2, Theorem 5]).

Theorem 14. Let $G$ be an open bounded symmetric subset of $H$ containing the origin and let $F$ be a multi-mapping $\bar{G} \rightarrow 2^{H}$ satisfying $F(-u)=-F(u)$ for all $u \in \partial G$. Then if $F \in \mathscr{F}_{G}(m Q M)$, the inclusion $0 \in F(u)$ is almost solvable in $\bar{G}$, i.e., $0 \in \overline{F(\bar{G})}$.

For the class $\mathscr{F}_{G}(m Q M)$ we have the following surjectivity theorem (see [2, Theorem 6]).

Theorem 15. Let $F: H \rightarrow 2^{H}$ be a mapping of the class $\mathscr{F}_{H}(m Q M)$ satisfying the following condition:
for any $K \in \mathbb{R}^{+}$there exists $R_{0} \in \mathbb{R}$ such that

$$
\begin{equation*}
\left|f^{*}\right|>K \quad \text { for all selections } f^{*} \in F(u),|u| \geqslant R_{0} . \tag{A}
\end{equation*}
$$

Assume moreover that there exists a positive real constant $r_{0}$ such that one of the following condition holds:

$$
\begin{gather*}
\frac{\left(f^{*}, u\right)}{|u|}+\left|f^{*}\right|>0 \quad \text { for all }|u| \geqslant r_{0} \text { and for all } f^{*} \in F(u)  \tag{B}\\
F(-u)=-F(u) \text { for all }|u| \geqslant r_{0}
\end{gather*}
$$

(C)

Then $\overline{F(H)}=H$.
Proof. Let $y \in H$. By (A), there exists $r_{1}>r_{0}$ such that $0 \notin \overline{F\left(\partial B_{r}\right)}$ for all $r \geqslant r_{1}$, where $B_{r}$ is the ball $B_{r}=\{u \in H| | u \mid<r\}$.

Assume first that (B) holds. It gives

$$
0 \notin(1-t) u+t F(u) \quad \text { for all } t \in[0,1] \text { and } u \in \partial B_{r} .
$$

By using the same argument as in the proof of Theorem 12 (with $R=I$ ) we can conclude (see (6)) that there exists a positive constant $\varepsilon_{0}=\varepsilon_{0}(r)$ such that

$$
0 \notin(1-t) u+t(F(u)+\varepsilon P u) \quad \text { for all }|u|=r, r>r_{1}, t \in[0,1] \text { and } 0<\varepsilon<\varepsilon_{0} .
$$

Hence

$$
d\left(F+\varepsilon P, B_{r}, 0\right)=+1 \quad \text { for all } r>r_{1} \text { and } 0<\varepsilon<\varepsilon_{0}
$$

We take $r_{2}>r_{1}$ such that $\left|f^{*}\right| \geqslant(3|y|+1)$ for all $f^{*} \in F(u),|u| \geqslant r_{2}$. Let $0<\varepsilon_{1}<\varepsilon_{0}$ be such that $\varepsilon_{1} r_{2}<|y|+1$. Consequently,

$$
\left|f^{*}+\varepsilon P u\right| \geqslant 3|y|-|y|=2|y|
$$

for all $f^{*} \in F(u), u \in \partial B_{r_{2}}$ and $0<\varepsilon<\varepsilon_{1}$. Since $|t y| \leqslant|y|$ we have

$$
t y \notin(F+\varepsilon P)\left(\partial B_{r_{2}}\right)
$$

and

$$
d\left(F+\varepsilon P, B_{r_{2}}, y\right)=d\left(F+\varepsilon P, B_{r_{2}}, 0\right)=+1 \quad \text { for all } 0<\varepsilon<\varepsilon_{1} .
$$

Hence $y \in(F+\varepsilon P)\left(B_{r_{2}}\right)$ for all $0<\varepsilon<\varepsilon_{1}$. It follows that $y \in \overline{F\left(B_{r_{2}}\right)} \subset \overline{F(H)}$.
Now we assume that (C) holds. We can use similar argument that in the first case. By applying Theorem 14, from (C) we obtain

$$
d\left(F+\varepsilon P, B_{r_{2}}, y\right)=d\left(F+\varepsilon P, B_{r_{2}}, 0\right)=\text { odd }
$$

for all $0<\varepsilon<\varepsilon_{1}$. Hence $y \in(F+\varepsilon P)\left(B_{r_{2}}\right)$ for all $0<\varepsilon<\varepsilon_{1}$ and the conclusion $y \in \overline{F(H)}$ follows. Thus the proof is completed.

It is not easy to check whether a multi-mapping $F$ is of the class $\mathscr{F}_{G}(m Q M)$ or $\mathscr{F}_{G}\left(m S_{+}\right)$. The simplest condition to check is the monotonicity and therefore the most useful theorems concern the class $\mathscr{F}_{G}(m P M)$. We prove the following existence result (see [2, Theorem 9]).

Theorem 16. Let $N \in(m P M)$, let $L: D(L) \subset H \rightarrow H$ be the linear mapping defined as above and let $c>0,-c \notin \sigma(L)$, where $\sigma(L)$ is the spectrum of $L$. If there exist real positive constants $\tau, \tau<\operatorname{dist}(-c, \sigma(L))$ and $K$ such that

$$
\left|f^{*}-c u\right| \leqslant(\tau|u|+K) \quad \text { for all } u \in H, f^{*} \in N(u)
$$

then for all $h \in H$ the inclusion (SInc) admits a solution in $H$, i.e. there exists $u \in D(L) \cap G$ such that $h \in L u+N(u)$.

Proof. We consider the homotopy

$$
F_{t}(u)=L u+c u+t(N(u)-c u) \quad \text { for all } t \in[0,1]
$$

and $y_{t}=t h$. Clearly $F_{0}(u)=L u+c u$ and $F_{1}(u)=L u+N(u)$. Since $-c \notin \sigma(L)$ the operator $(L+c I)^{-1}$ exists and is bounded. Thus we can put $r=\|(L+c I)\|^{-1}(K+$ $|h|+1) /\left(1-\tau\left\|(L+c I)^{-1}\right\|\right)$ and $G=B_{r}$. If for some $u \in \partial G \cap D(L)$ we have

$$
t h \in L u+c u+t(N(u)-c u)
$$

for some $t \in[0,1]$ then

$$
u=(L+c I)^{-1}\left(-t\left(f^{*}-c u\right)+t h\right), \quad f^{*} \in N(u)
$$

Hence

$$
r=|u| \leqslant\left\|(L+c I)^{-1}\right\|(\tau r+K+|h|)<r,
$$

a contradiction. Thus

$$
t h \notin L u+c u+t(N(u)-c u) \quad \text { for all } t \in[0,1] \text { and } u \in \partial G \cap D(L) .
$$

For $t=0$ we have $L u+c u=0$ not possessing a non-zero solution. The mapping $L+c I$ is linear, odd and $c I \in\left(m S_{+}\right)$, therefore using Lemma 3, Proposition 9 and then applying Theorem 13, we conclude $h \in F_{1}(G)$, i.e. $h \in L u+N(u)$ for some $u \in D(L) \cap G$.

## 6. $M$-REGULAR MULTI-FUNCTIONS

A function $p: \Omega \times \mathbb{R} \rightarrow \mathbb{R}$ is called superpositionally measurable [6] if $p(x, t, u(x, t))$ is measurable for any measurable function $u: \Omega \rightarrow \mathbb{R}$. A multi-function $S: \Omega \times$ $\mathbb{R} \rightarrow 2^{\mathbb{R}}$ is called measurable-bounded if there exist two superpositionally measurable functions $q_{-}(x, t, u)$ and $q_{+}(x, t, u)$ such that

$$
q_{-}(x, t, u) \leqslant q_{+}(x, t, u) \quad \text { and } \quad S(x, t, u)=\left[q_{-}(x, t, u), q_{+}(x, t, u)\right]
$$

for any $(x, t, u) \in \Omega \times \mathbb{R}$ where the function $q_{-}(x, t, u)$ is lsc in $u$, the function $q_{+}(x, t, u)$ is usc in $u$ and there exist positive constants $d_{1}, d_{2}$ and $c_{1}, c_{2} \in L^{2}(\Omega)$ such that

$$
\left|q_{-}(x, t, u)\right| \leqslant c_{1}(x, t)+d_{1}|u| \quad \text { and } \quad\left|q_{+}(x, t, u)\right| \leqslant c_{2}(x, t)+d_{2}|u|
$$

for any $(x, t, u) \in \Omega \times \mathbb{R}$. We denote by $(m M B)$ the set of all measurable-bounded multi-functions. By using a multi-function $S \in(m M B)$, for any $u \in L^{2}(\Omega)$ we put

$$
\begin{aligned}
N(u) & =\left\{v \in L^{2}(\Omega) \mid v(x, t) \in S(x, t, u(x, t))\right\} \\
& =\left\{v \in L^{2}(\Omega) \mid q_{-}(x, t, u(x, t)) \leqslant v(x, t) \leqslant q_{+}(x, t, u(x, t))\right\}
\end{aligned}
$$

and call it an $M$-regular multi-function. We denote the set of all such multi-functions by $(m M r) . N(u)$ is nonempty because $q_{ \pm}(x, t, u(x, t)) \in N(u)$. The degree theory we have constructed works just for multi-mappings that are w-usc.

Lemma 17. Let $N \in(m M r)$ and $u_{n} \rightarrow u$ in $L^{2}(\Omega)$. If a sequence $\left\{w_{n}^{*}\right\}$ satisfies $w_{n}^{*} \in N\left(u_{n}\right)$ and $w_{n}^{*} \rightharpoonup w^{*}$ in $L^{2}(\Omega)$ then $w^{*} \in N(u)$, i.e. $N: L^{2}(\Omega) \rightarrow 2^{L^{2}(\Omega)}$ is w-usc.

Proof. First, we have

$$
\begin{equation*}
q_{-}\left(x, t, u_{n}(x, t)\right) \leqslant w_{n}^{*}(x, t) \tag{11}
\end{equation*}
$$

for every $n \in \mathbb{N}$. Since $w_{n}^{*} \rightharpoonup w^{*}$, using the Mazur theorem we can choose a sequence $\left\{v_{n}\right\}, v_{n} \in \operatorname{conv}\left\{w_{n}^{*}, w_{n+1}^{*}, \ldots\right\}$ such that $v_{n} \rightarrow w^{*}$ almost everywhere in $L^{2}(\Omega)$. Thus we have

$$
v_{n}=\sum_{k=n}^{m_{n}} \lambda_{n, k} w_{k}^{*} ; \quad 0 \leqslant \lambda_{n, k} \leqslant 1 ; \quad \sum_{k=n}^{m_{n}} \lambda_{n, k}=1 \quad \text { for } n<m_{n} \in \mathbb{N}, n \leqslant k \leqslant m_{n}
$$

From (11) we have

$$
\sum_{k=n}^{m_{n}} \lambda_{n, k} q_{-}\left(x, t, u_{k}(x, t)\right) \leqslant \sum_{k=n}^{m_{n}} \lambda_{n, k} w_{k}^{*}(x, t)=v_{n}(x, t)
$$

By virtue of the convergence in the measure we can assume that $v_{n}(x, t) \rightarrow w^{*}(x, t)$ almost everywhere in $(x, t)$ [10]. Let $\left(x_{0}, t_{0}\right) \in \Omega$ be such an element. The mapping $s \rightarrow q_{-}\left(x_{0}, t_{0}, s\right)$ is lsc and so for every $\varepsilon>0$ there exists a positive integer $n_{0}$ such that for every $k \geqslant n_{0}$ we have

$$
q_{-}\left(x_{0}, t_{0}, u\left(x_{0}, t_{0}\right)\right)-\varepsilon \leqslant q_{-}\left(x_{0}, t_{0}, u_{k}\left(x_{0}, t_{0}\right)\right) .
$$

Summing this inequality for $k=n, n+1, \ldots, m_{n}$ with weights $\lambda_{n, k}$ we get

$$
\begin{aligned}
\sum_{k=n}^{m_{n}} \lambda_{n, k}\left(q_{-}\left(x_{0}, t_{0}, u\left(x_{0}, t_{0}\right)\right)-\varepsilon\right) & \leqslant \sum_{k=n}^{m_{n}} \lambda_{n, k} q_{-}\left(x_{0}, t_{0}, u_{k}\left(x_{0}, t_{0}\right)\right) \\
q_{-}\left(x_{0}, t_{0}, u\left(x_{0}, t_{0}\right)\right)-\varepsilon & \leqslant v_{n}\left(x_{0}, t_{0}\right)
\end{aligned}
$$

for all $n \geqslant n_{0}$. Hence by the convergence $v_{n}\left(x_{0}, t_{0}\right) \rightarrow w^{*}\left(x_{0}, t_{0}\right)$ we have

$$
q_{-}\left(x_{0}, t_{0}, u\left(x_{0}, t_{0}\right)\right)-\varepsilon \leqslant w^{*}\left(x_{0}, t_{0}\right) \text { for every } \varepsilon>0 .
$$

Finally, we get

$$
q_{-}\left(x_{0}, t_{0}, u\left(x_{0}, t_{0}\right)\right) \leqslant w^{*}\left(x_{0}, t_{0}\right)
$$

as we need. Similar argument leads to

$$
w^{*}\left(x_{0}, t_{0}\right) \leqslant q_{+}\left(x_{0}, t_{0}, u\left(x_{0}, t_{0}\right)\right) .
$$

Thus $w^{*} \in N(u)$.

## 7. SEmilinear wave equations

We show how the previous results can be applied to the semilinear wave equation (Prob). We state the precise setting of (Prob) by putting

$$
\begin{array}{cl}
q_{-}(x, t, u)=g_{-}(u)+f(x, t, u), & q_{+}(x, t, u)=g_{+}(u)+f(x, t, u), \\
g_{+}(u)=\limsup _{s \rightarrow u} g(u), & g_{-}(u)=\liminf _{s \rightarrow u} g(u) .
\end{array}
$$

We note that $g_{ \pm}$are Borel measurable. By Lemma 17, the Nemytskij operator $N: H \rightarrow 2^{H}, H=L^{2}(\Omega)$ defined by

$$
N(u)=\left\{v \in L^{2}(\Omega) \mid q_{-}(x, t, u(x, t)) \leqslant v(x, t) \leqslant q_{+}(x, t, u(x, t))\right\}
$$

is bounded and w-usc. $N \in(m M O N)$ and hence by Proposition $1, N \in(m P M)$. Let $C^{2}$ be the set of twice continuously differentiable functions $v:[0, \pi] \times \mathbb{R} \rightarrow \mathbb{R}$ satisfying $v(0, \cdot)=v(\pi, \cdot)=0$ and $2 \pi$-periodic in $t \in \mathbb{R}$.

A weak $2 \pi$-periodic solution of (Prob) for $h \in H$ is any $u \in H$ satisfying
(wProb) $\left(u, v_{t t}-v_{x x}\right)+\left(u^{*}, v\right)=(h, v) \quad$ for some $u^{*} \in N(u)$ and for all $v \in C^{2}$.
Let $\varphi_{m, n}(x, t)=\pi^{-1} e^{i m t} \sin n x$ for all $m \in \mathbb{Z}, n \in \mathbb{Z}^{+}$. Each $u \in L^{2}(\Omega)$ has a representation

$$
u=\sum_{m \in \mathbb{Z}, n \in \mathbb{Z}^{+}}\left(n^{2}-m^{2}\right) u_{m, n} \varphi_{m, n},
$$

where $u_{m, n}=\left(u, \varphi_{m, n}\right)$ and $\bar{u}_{m, n}=u_{m, n}$, since $u$ is a real function. The abstract realization of the wave operator $\frac{\partial^{2}}{\partial t^{2}}-\frac{\partial^{2}}{\partial x^{2}}$ in $L^{2}(\Omega)$ is the linear operator $L: D(L) \rightarrow$ $L^{2}(\Omega)$ defined by

$$
L u=\sum_{m \in \mathbb{Z}, n \in \mathbb{Z}^{+}}\left(n^{2}-m^{2}\right) u_{m, n} \varphi_{m, n}
$$

where

$$
D(L)=\left\{u \in L^{2}(\Omega)\left|\sum_{m \in \mathbb{Z}, n \in \mathbb{Z}^{+}}\right| n^{2}-\left.m^{2}\right|^{2}\left|u_{m, n}\right|^{2}<\infty\right\}
$$

It can be shown that $u \in L^{2}(\Omega)$ is a weak solution of (wProb) if and only if $h \in L u+N(u)$ with $u \in D(L)$. Moreover, $L$ is densely defined, self-adjoint, closed, $\operatorname{Im} L=(\operatorname{Ker} L)^{\perp}$ and $L$ has a pure point spectrum of eigenvalues $\sigma(L)=\left\{n^{2}-m^{2} \mid\right.$ $\left.m \in \mathbb{Z}, n \in \mathbb{Z}^{+}\right\}$with the corresponding eigenfunctions $\left\{\varphi_{m, n}\right\}$. Clearly $\sigma(L)$ is unbounded both from above and from below, any eigenvalue $\lambda \neq 0$ has a finite multiplicity, but Ker $L$ is infinite dimensional. The right inverse of $L_{0}, L_{0}^{-1}: \operatorname{Im} L \rightarrow \operatorname{Im} L$ is compact. Hence (Prob) fulfils all conditions for the inclusion (SInc) of Lemma 3. Theorem 16 immediately implies the following result (see [2, pp. 961-962]).

Theorem 18. The semilinear wave equation (Prob) under the conditions mentioned above admits a weak $2 \pi$-periodic solution for any $h \in L^{2}(\Omega)$ provided that in addition

$$
a \leqslant f(x, t, u) / u \leqslant b \quad \forall(x, t) \in \Omega \text { and } \forall u \in \mathbb{R},|u|>R
$$

for positive constants $a, b, R$ such that $m^{2}-n^{2} \notin[a, b]$ for all $n \in \mathbb{Z}^{+}$and $m \in \mathbb{Z}$.

We leave to the reader other extensions of results on pp. 960-962 from [2] to (Prob). In the conclusion, we consider the nonlinear discontinuous vibrating string equation

$$
\begin{gathered}
u_{t t}-u_{x x}+g(u)=f(x, t) \\
u(0, \cdot)=u(\pi, \cdot)=0
\end{gathered}
$$

where $f \in L^{\infty}(\Omega), g: \mathbb{R} \rightarrow \mathbb{R}$ is nondecreasing and bounded on $\mathbb{R}$. This equation is considered in the form

$$
\begin{gather*}
u_{x x}-u_{t t}+f(x, t) \in\left[g_{-}(u), g_{+}(u)\right]  \tag{12}\\
u(0, \cdot)=u(\pi, \cdot)=0 .
\end{gather*}
$$

We put $g(+\infty)=\sup _{\mathbb{R}} g, g(-\infty)=\inf _{\mathbb{R}} g$. We extend Theorem 1 of $[4]$ to (12).
Theorem 19. Let $f$ have a decomposition

$$
f(x, t)=f_{1}(x, t)+f_{2}(x, t), \quad f_{1} \in(\operatorname{ker} L)^{\perp}, \quad f_{2} \in \operatorname{ker} L
$$

together with

$$
g(-\infty)+\delta \leqslant f_{2}(x, t) \leqslant g(+\infty)-\delta \quad \forall(x, t) \in \Omega
$$

for some $\delta>0$. Then (12) has a weak $2 \pi$-periodic solution.
Proof. For $1>\varepsilon>0$ fixed, we consider the problem

$$
\begin{gather*}
u_{x x}-u_{t t}+f(x, t) \in\left[g_{-}(u), g_{+}(u)\right]+\varepsilon u  \tag{13}\\
u(0, \cdot)=u(\pi, \cdot)=0
\end{gather*}
$$

Since $1>\varepsilon>0$ and $g$ is nondecreasing and bounded on $\mathbb{R}$, Theorem 16 implies the existence of a weak $2 \pi$-periodic solution $u$ of (13). Now we show that $u$ is bounded in $L^{\infty}(\Omega)$ uniformly for $\varepsilon>0$ small. From (13) we have

$$
\begin{equation*}
-L u+f=z+\varepsilon u, \quad z \in\left[g_{-}(u), g_{+}(u)\right] \tag{14}
\end{equation*}
$$

We decompose $u=u_{1}+u_{2}, z=z_{1}+z_{2}, u_{1}, z_{1} \in(\operatorname{ker} L)^{\perp}, u_{2}, z_{2} \in \operatorname{ker} L$. Then (14) implies

$$
\begin{align*}
& L u_{1}+\varepsilon u_{1}+z_{1}=f_{1}  \tag{15}\\
& f_{2}=z_{2}+\varepsilon u_{2} . \tag{16}
\end{align*}
$$

Since $f_{1}, z \in L^{\infty}(\Omega)$ for $\varepsilon>0$ small, we have $\left|u_{1}\right|_{L^{\infty}}<\tilde{C}$ (see [4, p. 417]). In what follows $\tilde{C}$ denotes positive constants independent of $\varepsilon$. By (16) also

$$
\begin{equation*}
\int_{0}^{2 \pi} \int_{0}^{\pi} f_{2} u_{2} \mathrm{~d} x \mathrm{~d} t=\int_{0}^{2 \pi} \int_{0}^{\pi} z_{2} u_{2} \mathrm{~d} x \mathrm{~d} t+\varepsilon \int_{0}^{2 \pi} \int_{0}^{\pi} u_{2}^{2} \mathrm{~d} x \mathrm{~d} t \geqslant \int_{0}^{2 \pi} \int_{0}^{\pi} z_{2} u_{2} \mathrm{~d} x \mathrm{~d} t . \tag{17}
\end{equation*}
$$

We take $M>0$ such that

$$
g(v) \geqslant g(+\infty)-\frac{\delta}{2} \quad \forall v \geqslant M, \quad g(v) \leqslant g(-\infty)+\frac{\delta}{2} \quad \forall v \leqslant-M
$$

Let us define

$$
\Sigma_{+}=\left\{(x, t) \mid u_{2}(x, t) \geqslant M+\tilde{C}\right\}, \quad \Sigma_{-}=\left\{(x, t) \mid u_{2}(x, t) \leqslant-M-\tilde{C}\right\}
$$

Then we have

$$
z(x, t) \geqslant g(+\infty)-\frac{\delta}{2} \quad \forall(x, t) \in \Sigma_{+}, \quad z(x, t) \leqslant g(-\infty)+\frac{\delta}{2} \quad \forall(x, t) \in \Sigma_{-} .
$$

Now (17) gives

$$
\begin{aligned}
& (g(+\infty)-\delta) \int_{\Sigma_{+}} u_{2} \mathrm{~d} x \mathrm{~d} t+(g(-\infty)+\delta) \int_{\Sigma_{-}} u_{2} \mathrm{~d} x \mathrm{~d} t+\underset{\left(\Sigma_{+} \cup \Sigma_{-}\right)^{\prime}}{\int} f_{2} u_{2} \mathrm{~d} x \mathrm{~d} t \\
& \geqslant \int_{0}^{2 \pi} \int_{0}^{\pi} f_{2} u_{2} \mathrm{~d} x \mathrm{~d} t \\
& \geqslant \iint_{\left(\Sigma_{+} \cup \Sigma_{-}\right)^{\prime}} z_{2} u_{2} \mathrm{~d} x \mathrm{~d} t+\left(g(+\infty)-\frac{\delta}{2}\right) \int_{\Sigma_{+}} u_{2} \mathrm{~d} x \mathrm{~d} t+\left(g(-\infty)+\frac{\delta}{2}\right) \int_{\Sigma_{-}} u_{2} \mathrm{~d} x \mathrm{~d} t,
\end{aligned}
$$

hence we arrive at

$$
\int_{\Sigma_{+} \cup \Sigma_{-}}\left|u_{2}\right| \mathrm{d} x \mathrm{~d} t \leqslant \tilde{C}
$$

Consequently, we obtain that $\left|u_{2}\right|_{L^{1}}<\tilde{C}$. Since $\operatorname{dim} \operatorname{ker} L=\infty$, this estimate is not enough. Now following the same argument like on p. 419 of [4], we obtain that $\left|u_{2}\right|_{L^{\infty}}<\tilde{C}$. Hence for any $\varepsilon>0$ small, any solution of (14) satisfies $|u|_{L^{\infty}}<\tilde{C}$. Then clearly $|u|_{L^{2}}<\tilde{C}$. By passing to the limit $\varepsilon \rightarrow 0_{+}$in (14) like for Theorem 13 , we obtain a weak $2 \pi$-periodic solution of (12).

We remark according to Theorem 19: If $-g(-\infty)=g(+\infty)>0$ then the condition

$$
\left|\int_{0}^{2 \pi} \int_{0}^{\pi} f_{2}(x, t) z(x, t) \mathrm{d} x \mathrm{~d} t\right| \leqslant g(+\infty) \int_{0}^{2 \pi} \int_{0}^{\pi}|z(x, t)| \mathrm{d} x \mathrm{~d} t, \quad \forall z \in \operatorname{ker} L
$$

is necessary, and the condition

$$
\left|\int_{0}^{2 \pi} \int_{0}^{\pi} f_{2}(x, t) z(x, t) \mathrm{d} x \mathrm{~d} t\right| \leqslant(g(+\infty)-\delta) \int_{0}^{2 \pi} \int_{0}^{\pi}|z(x, t)| \mathrm{d} x \mathrm{~d} t, \quad \forall z \in \operatorname{ker} L
$$

for some $0<\delta<g(+\infty)$ is sufficient for the weak $2 \pi$-periodic solvability of (12). Finally, we note that similar results hold for (Prob) with nonincreasing $f$ and $g$ in $u$.

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