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> THREE METHODS FOR THE STUDY OF SEMILINEAR EQUATIONS AT RESONANCE

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Three methods for the study of the solvability of semilinear equations with noninvertible linear parts are compared: the alternative method, the continuation method of Mawhin and a new perturbation method [22]-[27]. Some extension of the last method and applications to differential equations in Banach spaces are presented.

1. Introduction. Most of nonlinear differential, integral or, more generally, functional equations have the form

$$
\begin{equation*}
L x=N(x) \tag{1.1}
\end{equation*}
$$

(called semilinear) where $L$ is a linear and $N$ a nonlinear operator, in appropriate function spaces. Usually, $L$ is defined on a dense linear subspace $Y$ of a Banach space $X$, takes values in a Banach space $Z$ and is a closed operator, and $N: X \rightarrow Z$ is continuous. Moreover, if $L$ has a trivial null space ker $L$, then $L$ is surjective; if $\operatorname{ker} L$ is nontrivial, it has a finite dimension equal to the codimension of the range space $L(Y) \subset Z$ (such operators are called Fredholm of index 0 ). We do not study the case $\operatorname{ker} L=\{0\}-$ it reduces to the fixed point problem for $L^{-1} N$. The mapping $N$ is usually compact (or contractive or monotone or A-proper or ...) and a suitable topological degree theory works and gives a large number of results. The case $\operatorname{ker} L \neq\{0\}$ is more complicated; we then say that equation (1.1) is at resonance. In this survey, we are interested only in the existence of solutions to equation (1.1) at resonance. It should be noticed, however, that there are results concerning the number of solutions (see e.g. [2], [28]) or localization of at least one solution in a given set (it is usually the cone of nonnegative functions - see [25], [29]).

We present a few examples of resonance problems. The most typical are boundary value problems for second order ordinary differential equations:

$$
\begin{equation*}
x^{\prime \prime}+m^{2} x=f\left(t, x, x^{\prime}\right), \quad x(0)=x(\pi)=0 \tag{1.2}
\end{equation*}
$$

[^0]where $m=1,2, \ldots$ and $f:[0, \pi] \times \mathbb{R}^{2} \rightarrow \mathbb{R}($ see $[7],[13],[21])$ or
\[

$$
\begin{equation*}
x^{\prime}=f(t, x), \quad x(0)=x(1), \tag{1.3}
\end{equation*}
$$

\]

where $f:[0,1] \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ ([4] and references therein). Ordinary differential equations can be replaced by partial differential ones ([17], [31], [15]):

$$
\begin{equation*}
P x-\lambda_{0} x=f(t, x), \quad x \mid \partial \Omega=0 \tag{1.4}
\end{equation*}
$$

where $P$ is a linear strongly elliptic differential operator of the second order on $\Omega$ with smooth coefficients, $\lambda_{0}$ is its first eigenvalue and $f: \mathbb{R} \times \operatorname{cl} \Omega \rightarrow \mathbb{R}$, or (see [3], [6])
(1.5) $\quad u_{t t}-\Delta u=f(t, z, u), \quad u \mid \partial \Omega=0, \quad u 2 \pi$-periodic in $t$,
where $\Delta$ is the Laplace operator, $f:[0,2 \pi] \times \operatorname{cl} \Omega \times \mathbb{R} \rightarrow \mathbb{R}$ (here, $\operatorname{dim} \operatorname{ker} L=$ $\infty)$. Equations with retarded argument or integro-differential equations can also be considered within this framework; usually the retardation or the integral are included in the nonlinear term:

$$
\begin{equation*}
x^{\prime}=f\left(t, x, x_{h}\right), \quad x(0)=x(1) \tag{1.6}
\end{equation*}
$$

where $x_{h}(t)=x(h(t)), h$ is a measurable function with $h(t) \leq t$ (see [1]), or

$$
\begin{equation*}
x^{\prime \prime}+x=\int_{0}^{\pi} K(t, s) f(s, x(s)) d s, \quad x(0)=x(\pi)=0 \tag{1.7}
\end{equation*}
$$

where $K:[0, \pi]^{2} \rightarrow \mathbb{R}, f:[0, \pi] \times \mathbb{R} \rightarrow \mathbb{R}$. We see that the nonlinear operator $N$ is a superposition operator or its composition with linear ones, so we have two possible ways for the choice of suitable Banach spaces. If $f$ is continuous, then $Z$ is the space of continuous functions with the sup-norm

$$
\|z\|=\sup _{t}|z(t)|
$$

$X=Z$ or $X$ is the space of $C^{1}$-functions if $f$ also depends on the first derivative. Then the solutions of all problems are classical in the sense that the derivatives exist everywhere. The second possibility: $f$ is a Carathéodory function, i.e. $f(\cdot, x)$ is measurable for all $x, f(t, \cdot)$ is continuous for a.e. $t$, and $f$ satisfies some growth condition

$$
|f(t, x)| \leq a|x|^{\varrho}+b(t)
$$

with $a \geq 0, \varrho \in(0,1]$ and $b \in L^{r}$. These assumptions ensure that

$$
N(x)(t)=f(t, x(t))
$$

maps $X=L^{p}$ into $Z=L^{q}$ with suitable $p, q, r, \varrho$ (see [10]). If $f$ depends on derivatives, these variables should behave similarly to $x$ and the operator $N$ maps an appropriate Sobolev space into $L^{q}$. The solutions obtained are called strong, i.e. the differential equation is satisfied a.e. The above examples of boundary value problems have a common feature: they have
the form $B x=0$ where $B$ is a linear operator. This enables us to restrict the subspace $Y$ and $X$ by the condition $B x=0$. For example, if $f$ is continuous in (1.2), then $X=\left\{x \in C^{1}([0, \pi]): x(0)=x(\pi)=0\right\}, Z=C([0, \pi])$ and $Y=\left\{x \in X: x^{\prime \prime} \in C([0, \pi])\right\}$; if $f$ is a Carathéodory function then $X=H_{0}^{1}, Y=X \cap H^{2}$ and $Z=L^{2}$, where $H^{m}$ is the Sobolev space of all functions whose derivatives up to order $m$ are square integrable and $H_{0}^{m}$ is the closure in $H^{m}$ of the set of smooth functions with compact support contained in $(0, \pi)$.

One can also study the boundary conditions of the form $B x=B_{\mathrm{n}}(x)$ where $B_{\mathrm{n}}$ is a nonlinear continuous operator taking values in the same space as $B$. Hence one can consider

$$
\begin{equation*}
x^{\prime \prime}+m^{2} x=f\left(t, x, x^{\prime}\right), \quad x(0)=r_{0}, \quad x(\pi)=r_{\pi} \tag{1.8}
\end{equation*}
$$

instead of (1.2),

$$
x^{\prime}=f(t, x), \quad x(0)-x(1)=\int_{0}^{1} x(s) d s
$$

instead of (1.3), and so on. This is possible since we can use product spaces for $X, Y$ and $Z$. For instance, the boundary value problem (1.8) can be treated as equation (1.1) with $X=C^{1}([0, \pi]), Y=C^{2}([0, \pi]), Z=C([0, \pi])$ $\times \mathbb{R}^{2}$, and

$$
L x=\left(x^{\prime \prime}, x(0), x(\pi)\right), \quad N(x)=\left(f\left(\cdot, x(\cdot), x^{\prime}(\cdot)\right), r_{0}, r_{\pi}\right)
$$

It is easy to introduce similar changes for all examples and both assumptions on $f$ (continuity or the Carathéodory conditions).

We return to the first example (1.2) -it was most extensively studiedand notice that, even for a bounded function $f$, a solution need not exist:

$$
x^{\prime \prime}+m^{2} x=f(t), \quad x(0)=x(\pi)=0
$$

The problem is selfadjoint with one-dimensional kernel spanned by $w(t)=$ $\sin m t$, so the solvability is equivalent to

$$
\begin{equation*}
\int_{0}^{\pi} w(t) f(t) d t=0 \tag{1.9}
\end{equation*}
$$

and, for $f=w$, this fails. In fact, equation (1.2) with $f(t, x)=g(x)+$ $f(t), g$ continuous and $\lim _{x \rightarrow \pm \infty} g(x)=g_{ \pm}$, has a solution if (1.9) holds and $g_{+} g_{-}<0$ (see below; this is a special case of the Landesman-Lazer condition). One can generalize this by taking $g_{ \pm}$not necessarily finite, the exact limits replaced by limsup and liminf and even (1.9) replaced by a weaker condition (see [9]), but the growth of $g$ cannot be too fast. If we take $g(x)=-(2 m+1) x, f(t)=\sin (m+1) t$, then (1.9) will be satisfied but
the corresponding problem

$$
x^{\prime \prime}+(m+1)^{2} x=\sin (m+1) t, \quad x(0)=x(\pi)=0
$$

has no solution. It is easy to see that the upper bound on the growth of $g$ is the distance between the eigenvalue $m^{2}$ which stands in (1.2) and the nearest eigenvalue $(m+1)^{2}$ or $(m-1)^{2}$. This situation is typical and we have a restriction on the constant

$$
\begin{equation*}
\gamma=\limsup _{\|x\| \rightarrow \infty}\|N(x)\| /\|x\| \tag{1.10}
\end{equation*}
$$

It is worth noticing, in connection with the above considerations, that one can generate many theorems by replacing equation (1.1) with $L x+A x=$ $N(x)+A x$ where $A$ is a linear operator with the same (or larger) domain as $L$ and $L+A$ is invertible, and by applying one of fixed point methods. However, the author does not refer to these theorems as resonance results.
2. Alternative method. This method seems to be the most general one for the study of equation (1.1). We have $L: X \supset Y \rightarrow Z, N: X \rightarrow$ $Z$ and we do not assume anything about $\operatorname{ker} L$ and $L(Y)$. The unique assumption is: $X$ and $Z$ admit decompositions

$$
X=X_{0} \oplus X_{1}, \quad Z=Z_{0} \oplus Z_{1}
$$

(topological direct sums) such that $\operatorname{ker} L \subset X_{0}$ and $Z_{1} \subset L(Y)$. Denote by $P: X \rightarrow X($ resp. $Q: Z \rightarrow Z)$ the projection onto $X_{0}\left(\right.$ resp. $\left.Z_{0}\right)$ along $X_{1}$ (resp. $Z_{1}$ ). Suppose that there exists a linear operator $H: Z_{1} \rightarrow X_{1}$ such that

$$
\begin{align*}
H(I-Q) L x & =(I-P) x & & \text { for } x \in Y  \tag{2.1}\\
L(H z) & =z & & \text { for } z \in Z_{1}  \tag{2.2}\\
Q(L x) & =L(P x) & & \text { for } x \in Y \tag{2.3}
\end{align*}
$$

If $Z_{1}=L(Y)$, then $\operatorname{ker} L=X_{0} \cap Y$ by (2.3), and conditions (2.1), (2.2) mean that $H$ is the inverse of $L \mid X_{1} \cap Y: X_{1} \cap Y \rightarrow Z_{1}$. However, this simplification fails in some examples.

We shall show that, under (2.1)-(2.3), equation (1.1) is equivalent to the system

$$
\begin{gather*}
x=P x+H(I-Q) N(x),  \tag{2.4}\\
Q(L x-N(x))=0 . \tag{2.5}
\end{gather*}
$$

In fact, we get (2.4) (resp. (2.5)) by applying $H(I-Q)$ (resp. $Q$ ) to (1.1) and using (2.1). Conversely, applying $L$ to (2.4) and using (2.2) and (2.3), we obtain

$$
L x=Q(L x)+(I-Q) N(x),
$$

so $(I-Q)(L x-N(x))=0$ and we have (1.1) by (2.5).

Equation (2.4) is usually called auxiliary and equation (2.5) the bifurcation equation. Generally, there are two possible ways to find a solution of the system (2.4)-(2.5). If we know that the auxiliary equation has a unique solution $S(\widehat{x})$ for any $\widehat{x}=P x \in X_{0}$ :

$$
S(\widehat{x})=\widehat{x}+H(I-Q) N(S(\widehat{x})),
$$

then the bifurcation equation will have the form

$$
\begin{equation*}
Q(L-N)(S(\widehat{x}))=0 \tag{2.6}
\end{equation*}
$$

The unique solvability is obtained if $H(I-Q) N$ is contractive or monotone and the last equation is defined and takes values in $X_{0}$, which is usually a finite-dimensional space. Thus equation (2.6) can be studied with the use of finite-dimensional methods, the Brouwer degree for example. We shall show below a sample of such arguments.

The second approach to the system (2.4)-(2.5) lies in the study of the $\operatorname{map} T: X \times X_{0} \rightarrow X \times X_{0}, T(x, \widehat{x})=(\widehat{x}+H(I-Q) N(x), \widehat{x}+Q(L-N)(x))$. Its fixed points (or rather their first components) are exactly solutions of the system examined. The typical method for finding such fixed points is the Leray-Schauder degree theory and its generalizations. We refer to the excellent survey by Cesari [4] where many examples of this and other types are given.

Example. Consider the typical problem

$$
\begin{equation*}
x^{\prime \prime}+x=f(t, x), \quad x(0)=x(\pi)=0, \tag{2.7}
\end{equation*}
$$

where $f:[0, \pi] \times \mathbb{R} \rightarrow \mathbb{R}$ is a Carathéodory function which is bounded:

$$
\begin{equation*}
|f(t, x)| \leq M, \quad \text { for } t \in[0, \pi], x \in \mathbb{R} \tag{2.8}
\end{equation*}
$$

and satisfies the Lipschitz condition with respect to $x$ :
(2.9) $\quad\left|f\left(t, x_{1}\right)-f\left(t, x_{2}\right)\right| \leq M\left|x_{1}-x_{2}\right|, \quad$ for $t \in[0, \pi], x_{1}, x_{2} \in \mathbb{R}$.

Since the system of functions

$$
w_{k}(t)=\sqrt{2 / \pi} \sin k t, \quad k=1,2, \ldots,
$$

is orthonormal and complete in $X=Z=L^{2}(0, \pi)$, we can take the orthogonal projectors

$$
P x=Q x=\sum_{k=1}^{m}\left(w_{k}, x\right) w_{k}
$$

where $(\cdot, \cdot)$ stands for the $L^{2}$-scalar product and $m$ will be chosen later. The null space of $L x=x^{\prime \prime}+x$ equals ker $L=\operatorname{Lin}\left\{w_{1}\right\} \subset X_{0}=\operatorname{Lin}\left\{w_{1}, \ldots, w_{m}\right\}$ and by the Hilbert-Schmidt theory, $H$ should be defined by the formula

$$
H z=\sum_{k=2}^{\infty}\left(k^{2}-1\right)^{-1}\left(w_{k}, z\right) w_{k}
$$

$\left\{w_{k}\right\}$ is a system of eigenfunction for $x^{\prime \prime}$ and the corresponding sequence of eigenvalues is $\left\{k^{2}\right\}$. Conditions (2.1)-(2.3) are easily verified, hence we can study the auxiliary equation. By (2.9), the right-hand side of (2.4) is a contraction for sufficiently large $m$ :

$$
\begin{aligned}
\|H(I-Q)(N(x)-N(y))\|^{2} & =\left\|\sum_{k=m+1}^{\infty}\left(k^{2}-1\right)^{-1}\left(N(x)-N(y), w_{k}\right) w_{k}\right\|^{2} \\
& =\sum_{k=m+1}^{\infty}\left(k^{2}-1\right)^{-1}\left|\left(N(x)-N(y), w_{k}\right)\right|^{2} \\
& \leq M^{2} \sum_{k=m+1}^{\infty}\left(k^{2}-1\right)^{-1}\|x-y\|^{2}
\end{aligned}
$$

On the other hand, the functions $H(I-Q) N(x)$ are bounded:

$$
\begin{equation*}
|H(I-Q) N(x)(t)| \leq \sum_{k=m+1}^{\infty}\left(k^{2}-1\right)^{-1}(2 \pi)^{-1} M \tag{2.10}
\end{equation*}
$$

and the last constant tends to 0 as $m \rightarrow \infty$. Hence, for any $\widehat{x}=\sum_{k=1}^{m} a_{k} w_{k}$, there exists a unique $S \hat{x}$ such that

$$
S \widehat{x}=\widehat{x}+H(I-Q) N(S \widehat{x})
$$

and we can pass to equation (2.6) which assumes the form

$$
\int_{0}^{\pi}(S \widehat{x})^{\prime \prime} w_{k}+\int_{0}^{\pi}(S \widehat{x}) w_{k}=\int_{0}^{\pi} f(t, S \widehat{x}(t)) w_{k}(t) d t, \quad k=1, \ldots, m
$$

or, after integration by parts,

$$
\left(1-k^{2}\right) a_{k}=\int_{0}^{\pi} f(t, S \widehat{x}(t)) w_{k}(t) d t, \quad k=1, \ldots, m
$$

where $a_{k}=\left(w_{k}, \widehat{x}\right)$. Let us introduce the mapping $F: \mathbb{R}^{m} \rightarrow \mathbb{R}^{m}, F=$ $\left(F_{1}, \ldots, F_{m}\right)$, where

$$
\begin{aligned}
F_{k}\left(a_{1}, \ldots, a_{m}\right)= & \left(k^{2}-1\right) a_{k} \\
& +\int_{0}^{\pi} f\left(t, \sum_{k=1}^{m} a_{k} w_{k}(t)+H(I-Q) N(S \widehat{x})(t)\right) w_{k}(t) d t .
\end{aligned}
$$

If $k \neq 1$, then the sign of $F_{k}$ is determined by the sign of $a_{k}$ for $\left|a_{k}\right| \geq$ $\left(k^{2}-1\right)^{-1} M \sqrt{2 / \pi}$. Set $r_{k}=2\left(k^{2}-1\right)^{-1} M \sqrt{2 / \pi}$ for $k \geq 2$. The situation for $k=1$ is different. Using (2.9), we can estimate

$$
\left|\sum_{k=2}^{m} a_{k} w_{k}(t)+H(I-Q) N(S \widehat{x})(t)\right| \leq C
$$

On the other hand, $a_{1} w_{1}(t) \rightarrow \pm \infty$ as $a_{1} \rightarrow \pm \infty$ uniformly on any set $[\delta, \pi-\delta]$. Hence if we assume that

$$
\begin{equation*}
x f(t, x) \geq 0 \tag{2.11}
\end{equation*}
$$

for sufficiently large $|x|$, or

$$
\begin{equation*}
x f(t, x) \leq 0 \tag{2.12}
\end{equation*}
$$

for large $|x|$, where strict inequalities hold on a set of a positive measure, then we can find $r_{1}>0$ such that $F_{1}$ has opposite signs on both sides $a_{1}= \pm r_{1}$ of the cube $U=\left(-r_{1}, r_{1}\right) \times \ldots \times\left(-r_{m}, r_{m}\right)$. The same was proved above for $F_{k}, k \geq 2$, and $a_{k}= \pm r_{k}$. Therefore the Brouwer degree [18] $\operatorname{deg}(F, U, 0)$ is defined and $F$ can be homotopically deformed to the antipodal map $\widehat{F}$, $\widehat{F}(-a)=-\widehat{F}(a)$. By the Borsuk Antipodensatz [18],

$$
\operatorname{deg}(F, U, 0) \equiv 1(\bmod 2)
$$

and the equation $F(a)=0$ has a solution in $U$. But this is a solution of the bifurcation equation. We have proved that problem (2.7) has a solution provided that $f$ satisfies (2.8), (2.9) and (2.11) or (2.12).

Remark. Condition (2.12) is exactly as in the recent paper by Iannacci and Nkashama [13], however, they omit the strong assumption of the boundedness of $f$. Similarly, we can prove by the alternative method that problem (2.7) has a solution if, instead of (2.12), we have

$$
\begin{equation*}
\int_{0}^{\pi} f_{+}(t) \sin t d t<0<\int_{0}^{\pi} f_{-}(t) \sin t d t \tag{2.13}
\end{equation*}
$$

where

$$
f_{+}(t)=\limsup _{x \rightarrow+\infty} f(t, x), \quad f_{-}(t)=\liminf _{x \rightarrow-\infty} f(t, x)
$$

This is exactly the Landesman-Lazer condition in the form of [7].
3. Coincidence degree. Let us go back to the alternative scheme and suppose that $X_{0}=\operatorname{ker} L, Z_{1}=L(Y)$ and both $X_{0}$ and $Z_{0}$ are linear subspaces of the same finite dimension. Then there exists an isomorphism (not unique) $J: Z_{0} \rightarrow X_{0}$. Conditions (2.1), (2.3) are trivially satisfied and, by (2.2), $H$ should be a right inverse of $L$. The bifurcation equation can be rewritten as

$$
\begin{equation*}
J Q N(x)=0 \tag{3.1}
\end{equation*}
$$

Obviously, the system of equations (2.4), (3.1) is equivalent to one equation

$$
\begin{equation*}
x=P x+J Q N(x)+H(I-Q) N(x) . \tag{3.2}
\end{equation*}
$$

The nonlinear operator $N$ is called $L$-compact if $Q N$ and $H(I-Q) N$ are compact operators. Usually, $H$ is compact (maps bounded sets into compact
ones) and $N$ maps bounded sets into bounded ones, which implies the $L$ compactness of $N$. The main theorem for $L$-compact mappings was proved by J. Mawhin [19].

Theorem 1. Let $\Omega$ be an open bounded subset in $X$ such that
(i) the equation $L x=\lambda N(x)$ has no solution for $\lambda \in(0,1]$ and $x \in$ $Y \cap \partial \Omega$,
(ii) the following Brouwer degree is defined and does not vanish:

$$
\operatorname{deg}(J Q N \mid \operatorname{ker} L, \Omega \cap \operatorname{ker} L, 0) \neq 0
$$

Then equation (1.1) has a solution in $\Omega$.
Proof. Let $h(x, \lambda)=(P+J Q N)(x)+\lambda H(I-Q) N(x)$ for $\lambda \in[0,1]$, $x \in \operatorname{cl} \Omega$. Since $N$ is $L$-compact, the homotopy $h$ is compact. Moreover, fixed points of $h$ are exactly solutions of $L x=\lambda N(x)$, so $h$ does not admit fixed points on $\partial \Omega$ (for $\lambda>0$, this follows from (i); for $\lambda=0$, it is a consequence of the fact that the Brouwer degree is defined). Hence the Leray-Schauder degree

$$
\operatorname{deg}_{\mathrm{LS}}(I-h(\cdot, \lambda), \Omega, 0)
$$

is independent of $\lambda$ (comp. [18], [6]). But, for $\lambda=0$, it reduces to the Brouwer degree (ii). Therefore, $\operatorname{deg}_{\text {LS }}(I-h(\cdot, 1), \Omega, 0) \neq 0$, hence $h(\cdot, 1)$ has a fixed point in $\Omega$ which is a solution to (1.1).

The Brouwer degree (ii) is called the coincidence degree of $N$ with respect to $L$. We refer to [11] for more information about this theory. Notice also that this theorem is closely related to Fučik's theorem [10] where only Hilbert spaces are considered.

We present one application of the theorem (comp. [20]). Consider the periodic boundary value problem

$$
\begin{equation*}
x^{\prime}=f(t, x), \quad x(0)=x(T), \tag{3.3}
\end{equation*}
$$

where $f: \mathbb{R} \times \mathbb{R}^{n} \rightarrow \mathbb{R}^{n}$ is a continuous function, $T$-periodic with respect to $t$. Here, $X=Z=C_{T}\left(\mathbb{R}, \mathbb{R}^{n}\right)$, the space of continuous $T$-periodic functions, $Y=C_{T}^{1}\left(\mathbb{R}, \mathbb{R}^{n}\right), L x=x^{\prime}, N(x)(t)=f(t, x(t))$, ker $L$ is the space of constant functions which is naturally isomorphic to $\mathbb{R}^{n}$, and $L(Y)$ is the subspace of continuous functions for which the integral over $[0, T]$ vanishes. Obviously, ker $L \oplus L(Y)=X=Z$ and we can take $P=Q$ given by

$$
P x=\int_{0}^{T} x(t) d t \quad \text { (constant function). }
$$

Moreover, $J=I$, and $H: Z \supset L(Y) \rightarrow L(Y) \subset X$ is defined by

$$
H z(t)=\int_{0}^{t} z(s) d s
$$

All introductory assumptions are satisfied and we can pass to (i) and (ii). It is easy to see that (i) is guaranteed by

$$
\begin{equation*}
f_{i}(t, x) \neq 0 \quad \text { for }\left|x_{i}\right|=M_{i}, i=1, \ldots, n \tag{3.4}
\end{equation*}
$$

where $M_{i}$ are constants and $f_{i}$ denotes the $i$ th coordinate of $f$. The choice of $\Omega$ is simple:

$$
x \in \Omega \subset C_{T}\left(\mathbb{R}, \mathbb{R}^{n}\right) \Leftrightarrow\left|x_{i}(t)\right|<M_{i}, i=1, \ldots, n, t \in \mathbb{R}
$$

and the proof of (i) is obvious.
The mapping $J Q N \mid \operatorname{ker} L$ is given by the formula

$$
F(a)=\int_{0}^{T} f(t, a) d t, \quad a \in \mathbb{R}^{n}
$$

and $\Omega \cap \operatorname{ker} L$ is the cube $\left(-M_{1}, M_{1}\right) \times \ldots \times\left(-M_{n}, M_{n}\right)$. Thus we have to strengthen condition (3.4):

$$
\begin{equation*}
f_{i}\left(t, x_{1}, \ldots, x_{i-1},-M_{i}, \ldots, x_{n}\right) f_{i}\left(t, x_{1}, \ldots, x_{i-1}, M_{i}, \ldots, x_{n}\right)<0 \tag{3.5}
\end{equation*}
$$

for any $\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n}$ such that $\left|x_{i}\right| \leq M_{i}$ for all $i=1, \ldots, n$. This implies that $F_{i}$ takes opposite signs on both sides $a_{i}= \pm M_{i}$ of the cube. The arguments as in the example from Section 2 show that $\operatorname{deg}(F, \operatorname{ker} L \cap \Omega, 0)$ is odd and, therefore, does not vanish.

Remark. One can study the same problem in another way:

$$
\frac{d}{d t}\|x\|^{2}=2\left(x^{\prime}, x\right)=2(f(t, x), x)
$$

If $(f(t, a), a) \neq 0$ for $\|a\|=R$, then $\Omega=B(0, R)$ satisfies (i). Moreover, the Krasnosel'skiĭ condition

$$
(F(a), a) \leq 0 \quad \text { for }\|a\|=R
$$

implies (ii) (comp. [18]). Thus the following condition is sufficient for the solvability of (3.3):

$$
(f(t, a), a)<0 \quad \text { for }\|a\|=R, t \in[0, T] .
$$

4. The perturbation method. The third method for the study of resonance problems is based on the observation that a linear noninvertible operator $L$ after an arbitrarily small perturbation, $L+\lambda I$ for instance, becomes invertible. One can solve perturbed equations and find conditions that ensure the existence of a convergent sequence of solutions. Its limit is the solution to (1.1) we look for. This rather natural observation was first made by P. Hess in his elegant proof of the Landesman-Lazer theorem [12]. The abstract result was given by de Figueiredo [5] and, then, more specific methods attracted the attention of mathematicians. Recently, the present author has applied this approach to resonance problems independently of the above
mentioned papers. The use of inverse operators of perturbed linear maps enables us to pass to unbounded nonlinearities and even to nonlinearities with linear growth. The method is developed in a series of papers [22]-[27] where also applications to boundary value problems are given.

Let $L(\lambda)$ be a continuous family of linear operators $Y \rightarrow Z$ for $\lambda$ from a nbhd of $\lambda_{0} \in \mathbb{R}$. We assume that $Y$ and $Z$ are Banach spaces such that all $L(\lambda)$ are bounded and the continuity of $L(\cdot)$ is meant to be the norm continuity. Let $L(\lambda)$ be invertible for $\lambda \neq \lambda_{0}$, and $L\left(\lambda_{0}\right)$ be a Fredholm linear operator (of index zero). Let the inclusion map $J: Y \rightarrow X$ be compact. Denote by $w_{1}, \ldots, w_{n}$ a basis of ker $L$ and suppose that the limits

$$
\lim _{\lambda \rightarrow \lambda_{0}}\left\|L(\lambda) w_{j}\right\|^{-1} L(\lambda) w_{j}, \quad j=1, \ldots, n
$$

exist, spanning a topological complement of $L\left(\lambda_{0}\right)(Y)$ (these limits always exist for $\left.L(\lambda)=L\left(\lambda_{0}\right)+\left(\lambda-\lambda_{0}\right) I\right)$. Under this assumption, the inverse operators $G(\lambda)=L(\lambda)^{-1}$ have the form

$$
G(\lambda)=G_{0}(\lambda)+\sum_{j=1}^{n} c_{j}(\lambda)\left\langle u_{j}(\lambda), \cdot\right\rangle w_{j}
$$

where $G_{0}(\lambda): Z \rightarrow Y, c_{j}(\lambda) \in \mathbb{R}, u_{j}(\lambda) \in Z^{*}$, while $G_{0}$ and $u_{j}$ have continuous extensions to $\lambda_{0}$ and

$$
\lim _{\lambda \rightarrow \lambda_{0}}\left|c_{j}(\lambda)\right|=\infty, \quad j=1, \ldots, n
$$

(comp. [26]). Moreover, $G_{0}(\lambda)$ takes values in a fixed space which is a topological complement of ker $L$. Notice that

$$
L\left(\lambda_{0}\right)(Y)=\bigcap_{j=1}^{n} \operatorname{ker} u_{j}\left(\lambda_{0}\right)
$$

The examined equation

$$
\begin{equation*}
L\left(\lambda_{0}\right) y=N(J y) \tag{4.1}
\end{equation*}
$$

is equivalent to the system

$$
\begin{gather*}
x=J G_{0}\left(\lambda_{0}\right) N(x)+\sum d_{j} J w_{j},  \tag{4.2}\\
\left\langle u_{j}\left(\lambda_{0}\right), N(x)\right\rangle=0, \quad j=1, \ldots, n, \tag{4.3}
\end{gather*}
$$

where $x=J y$ and the first summand in (4.2) is an element of a topological complement $\widetilde{X}$ of $\bar{X}=\operatorname{Lin}\left\{J w_{1}, \ldots, J w_{n}\right\}$. Fix $\lambda_{1}$ close to $\lambda_{0}$.

Theorem 2. Suppose there exist open bounded subsets $\widetilde{U}$ of $\widetilde{X}$ and $U$ of $\mathbb{R}^{n}$ such that
(*) there is no solution $\left(\widetilde{x},\left(d_{1}, \ldots, d_{n}\right)\right)$ to the system

$$
\widetilde{x}=\left(\lambda-\lambda_{1}\right)\left(\lambda_{0}-\lambda_{1}\right)^{-1} J G_{0}(\lambda) N(x),
$$

$$
d_{j}=c_{j}(\lambda)\left\langle u_{j}(\lambda), N(x)\right\rangle, \quad j=1, \ldots, n
$$

on the boundary of $\widetilde{U} \times U$ where $x=\widetilde{x}+\sum d_{j} J w_{j}$.
If we denote $g=\left(g_{1}, \ldots, g_{n}\right): \operatorname{cl} U \rightarrow \mathbb{R}^{n}$ by the formula

$$
g_{j}\left(d_{1}, \ldots, d_{n}\right)=d_{j}-c_{j}\left(\lambda_{1}\right)\left\langle u_{j}\left(\lambda_{1}\right), N\left(\sum d_{i} J w_{i}\right)\right\rangle
$$

and if

$$
\operatorname{deg}(g, U, 0) \neq 0
$$

then equation (4.1) has a solution $y$ such that $J y=\widetilde{x}+\sum d_{j} J w_{j}$ and $\widetilde{x} \in$ cl $\widetilde{U},\left(d_{1}, \ldots, d_{n}\right) \in \operatorname{cl} U$.

Proof. Define an open bounded set $V \subset X$ by $V=\left\{\widetilde{x}+\sum d_{j} J w_{j}:\right.$ $\left.\widetilde{x} \in \widetilde{U},\left(d_{1}, \ldots, d_{n}\right) \in U\right\}$. Then condition $(*)$ means that the homotopy

$$
H(x, \lambda)=\frac{\lambda-\lambda_{1}}{\lambda_{0}-\lambda_{1}} J G_{0}(\lambda) N(x)+\sum_{j=1}^{n} c_{j}(\lambda)\left\langle u_{j}(\lambda), N(x)\right\rangle J w_{j}
$$

for $x \in X$ and $\lambda \in\left(\lambda_{0}, \lambda_{1}\right]$, is fixed point free on the boundary $\partial V$. It follows that the Leray-Schauder degree

$$
\operatorname{deg}_{\mathrm{LS}}(I-H(\cdot, \lambda), V, 0)
$$

is defined and independent of $\lambda \in\left(\lambda_{0}, \lambda_{1}\right]$. On the other hand, $H\left(\cdot, \lambda_{1}\right)$ has a finite-dimensional range and, therefore,

$$
\operatorname{deg}_{\mathrm{LS}}\left(I-H\left(\cdot, \lambda_{1}\right), V, 0\right)=\operatorname{deg}(g, U, 0) \neq 0
$$

Thus, for any $\lambda>\lambda_{0}$, there exists a solution to the equation $x=H(x, \lambda)$ in $V$.

Take $\lambda_{k} \rightarrow \lambda_{0}$ and $x_{k}=\widetilde{x}_{k}+\sum_{j} d_{j}^{k} J w_{j}=H\left(x_{k}, \lambda_{k}\right)$ such that $\widetilde{x}_{k} \in \widetilde{U}$, $\left(d_{1}^{k}, \ldots, d_{n}^{k}\right) \in U$. Since $J$ is compact and $\left(x_{k}\right)$ is bounded, we can choose subsequences (still denoted by $\left(\lambda_{k}\right)$ and $\left(x_{k}\right)$ for simplicity) such that

$$
\begin{aligned}
\frac{\lambda_{k}-\lambda_{1}}{\lambda_{0}-\lambda_{1}} J G_{0}\left(\lambda_{k}\right) N\left(x_{k}\right) & \rightarrow x_{0} \\
c_{j}\left(\lambda_{k}\right)\left\langle u_{j}\left(\lambda_{k}\right), N\left(x_{k}\right)\right\rangle & \rightarrow d_{j}, \quad j=1, \ldots, n
\end{aligned}
$$

Hence $x_{k} \rightarrow x_{0}+\sum d_{j} J w_{j}=: x$. Obviously, $x$ satisfies (4.2)-(4.3).
This theorem is closely related to Theorem 1 . Since $\lambda_{1}$ is near $\lambda_{0}$, the real coefficients $c_{j}\left(\lambda_{1}\right)$ are arbitrarily large, so the degree of $g$ equals (up to sign) the degree of the map $\left(c_{j}\left(\lambda_{1}\right)^{-1} g_{j}\right)_{j \leq n}$, and the latter equals $\operatorname{deg}(h, U, 0)$ with

$$
h_{j}\left(d_{1}, \ldots, d_{n}\right)=\left\langle u_{j}\left(\lambda_{0}\right), N\left(\sum d_{i} J w_{i}\right)\right\rangle .
$$

The last mapping is exactly $J Q N \mid \operatorname{ker} L$ from Theorem 1 but here it is described better for applications. It should also be noticed that (4.2) is the
auxiliary equation and (4.3) the bifurcation equation from the alternative method.

Theorem 2 produces results concerning nonlinearities with restricted growth. If $N$ is sublinear:

$$
\lim _{\|x\| \rightarrow \infty}\|N(x)\| /\|x\|=0
$$

then the condition guaranteeing the solvability of equation (4.1) has the form: for any $\left(x_{\nu}\right) \subset X$ with $\left\|x_{\nu}\right\| \rightarrow \infty$ and $\left\|x_{\nu}\right\|^{-1} x_{\nu} \rightarrow \sum d_{j} J w_{j}$, there exists $j_{0} \in\{1, \ldots, n\}$ such that

$$
\limsup _{\nu \rightarrow \infty} d_{j_{0}}\left\langle u_{j_{0}}\left(\lambda_{0}\right), N\left(x_{\nu}\right)\right\rangle<0
$$

We refer to it as L-L. The same is true if we reverse the inequality and replace limsup by liminf. The proof via Theorem 2 can be found in [24], and the direct one in [23]. The case of nonlinearities with linear growth is more difficult. If

$$
\gamma=\lim _{\|x\| \rightarrow \infty}\|N(x)\| /\|x\| \in(0, \infty)
$$

then we have a natural restriction on $\gamma$ :

$$
\gamma\left\|J G_{0}\left(\lambda_{0}\right)\right\|<1
$$

and the L-L condition has to be replaced by a stronger one, $\mathrm{L}-\mathrm{L}_{1}$ : there exist $R>0$ and $\sigma>\gamma\left\|J G_{0}\left(\lambda_{0}\right)\right\|\left(1-\gamma\left\|J G_{0}\left(\lambda_{0}\right)\right\|\right)^{-1}$ such that, for any $j=1, \ldots, n$ and $\left|d_{j}\right|>R$,

$$
\begin{aligned}
& \sup \left\{d_{j}\left\langle u_{j}\left(\lambda_{0}\right), N\left(\widetilde{x}+\sum d_{i} J w_{i}\right)\right\rangle:\right. \\
& \left.\qquad\left|d_{i}\right| \leq\left|d_{j}\right|, \widetilde{x} \in \widetilde{X},\|\widetilde{x}\| \leq \sigma\left\|\sum d_{i} J w_{i}\right\|\right\}<0
\end{aligned}
$$

We can replace the last inequality by inf $>0$, but for all $j$ simultaneously. Here, the proof can be carried out only via homotopy arguments, as in Theorem 2 of [24]. That paper also contains a result on superlinear nonlinearities ( $N$ is called superlinear if $\|N(x)\| \leq a+b\|x\|^{\varrho}$ with $\varrho>1$ ).

The abstract theorems have numerous applications to boundary value problems. The first application found was the existence result for partial differential equations, obtained via the alternative method by Landesman and Lazer [17]. We present a slight generalization of it but refer to [23] where the full discussion is given.

Consider the Dirichlet boundary value problem

$$
\begin{equation*}
P u-\lambda_{0} u=f(x, u), \quad u \in H_{0}^{m}(\Omega), \tag{4.4}
\end{equation*}
$$

where $P$ is a selfadjoint strongly elliptic operator of order $2 m$ on an open bounded domain $\Omega \subset \mathbb{R}^{k}$ with Lipschitzian boundary, $\lambda_{0}$ is an eigenvalue for $P$, i.e. $P u-\lambda_{0} u=0$ has nontrivial solutions belonging to $H_{0}^{m}(\Omega)$,
$f: \Omega \times \mathbb{R} \rightarrow \mathbb{R}$ is a Carathéodory function (measurable with respect to $x$ and continuous with respect to $u$ ) with a sublinearly restricted growth:

$$
\begin{equation*}
\lim _{|u| \rightarrow \infty} \underset{x}{\operatorname{esssup}}|f(x, u)| /|u|=0 . \tag{4.5}
\end{equation*}
$$

Since $f$ need not be continuous, we should look for strong solutions of the problem. Put $X=Z=L^{2}(\Omega)$ and $Y=H^{2 m}(\Omega) \cap H_{0}^{m}(\Omega)$ where $H^{m}(\Omega)$ is the Sobolev space of all functions whose distributional derivatives up to order $m$ are in $L^{2}(\Omega)$ and $H_{0}^{m}(\Omega)$ is the closure in this space of the set of all smooth functions with supports contained in $\Omega$. The condition $u \in H_{0}^{m}(\Omega)$ corresponds to the null Dirichlet data $D^{\alpha} u \mid \partial \Omega=0$ for $|\alpha|<m$.

Let $L(\lambda)=P-\lambda I$ and $N(u)(x)=f(x, u(x))$. By the Hilbert-Schmidt theory,

$$
G(\lambda)=\sum_{j=n+1}^{\infty}\left(\lambda_{j}-\lambda\right)^{-1}\left(w_{j}, \cdot\right) w_{j}+\sum_{j=1}^{n}\left(\lambda_{j}-\lambda\right)^{-1}\left(w_{j}, \cdot\right) w_{j}
$$

where $\left\{\lambda_{j}: j=1,2, \ldots\right\}$ is the sequence of all eigenvalues of $P,\left\{w_{j}\right.$ : $j=1,2, \ldots\}$ is the sequence of the corresponding eigenfunctions which form an orthonormal basis in $L^{2}(\Omega),(\cdot, \cdot)$ is the $L^{2}$-scalar product. We choose $\lambda_{1}=\ldots=\lambda_{n} \neq \lambda_{j}$ for $j>n$ and denote this eigenvalue by $\lambda_{0}$. The first summand (denote it by $G_{0}(\lambda)$ ) has a continuous extension to $\lambda=\lambda_{0}$. Moreover,

$$
c_{j}(\lambda)=\left(\lambda_{0}-\lambda\right)^{-1}, \quad\left\langle u_{j}(\lambda), z\right\rangle=\left(w_{j}, z\right),
$$

in our notation. The inclusion map $J: H^{2 m} \cap H_{0}^{m} \rightarrow L^{2}$ is obviously compact. By (4.5), $N$ is sublinear and we can use the L-L condition. It now has the form: for any $\left(u_{\nu}\right) \subset L^{2}$ such that $\left\|u_{\nu}\right\| \rightarrow \infty$ and $\left\|u_{\nu}\right\|^{-1} u_{\nu} \rightarrow \sum d_{j} w_{j}$, there exists $j_{0}$ such that

$$
\limsup _{\nu \rightarrow \infty} d_{j_{0}} \int_{\Omega} w_{j_{0}}(x) f\left(x, u_{\nu}(x)\right) d x<0
$$

or $\lim \inf \ldots>0$. If we add all inequalities, we shall get a stronger assumption but having some advantages. Suppose that the limits

$$
\begin{equation*}
\lim _{u \rightarrow \pm \infty} f(x, u)=f_{ \pm}(x) \tag{4.6}
\end{equation*}
$$

exist. Then the following condition is sufficient for the solvability of (4.4): for any $d \in \mathbb{R}^{n}$ such that $\|d\|=1$, the numbers

$$
\begin{aligned}
& \int_{A_{+}} f_{+}(x) \sum d_{j} w_{j}(x) d x+\int_{A_{-}} f_{-}(x) \sum d_{j} w_{j}(x) d x \\
& \int_{A_{+}} f_{-}(x) \sum d_{j} w_{j}(x) d x+\int_{A_{-}} f_{+}(x) \sum d_{j} w_{j}(x) d x
\end{aligned}
$$

have opposite signs, where $A_{+}=\left\{x: \sum d_{j} w_{j}(x)>0\right\}$ and $A_{-}=\{x:$ $\left.\sum d_{j} w_{j}(x)<0\right\}$ (cf. [31]). The one-dimensional resonance ( $n=1$ ) was studied by Landesman and Lazer [17]: the sufficient condition is then simpler since we can take only $d_{1}=+1$ (for $d_{1}=-1$, we get reversed signs).

The limits in (4.6) can be replaced by limsup and liminf and the nonlinearity $f$ can also depend also on derivatives (comp. [23]). Moreover, the method can be applied to all equations with selfadjoint operators having inverses compact in $L^{2}$. This enables us to study periodic problems via the perturbation method [25].

If the nonlinearity $f$ in problem (4.4) has a linear growth:

$$
\limsup _{|u| \rightarrow \infty} \underset{x}{\operatorname{esssup}}|f(x, u)| /|u|=\gamma \in(0, \infty),
$$

we have no integral conditions of the above type that ensure the existence of solutions to (4.4). We have only found the following sufficient condition if $n=1$ [23]:

$$
\begin{equation*}
\check{a} \leq \frac{f(x, u)-b(x)}{u} \leq \hat{a} \quad \text { for } x \in \Omega,|u|>M, \tag{4.7}
\end{equation*}
$$

where $\check{a}, \hat{a}, M$ are positive constants, $b \in L^{2}(\Omega)$ and

$$
\gamma\left(\min _{j \neq 0}\left|\lambda_{0}-\lambda_{j}\right|\right)^{-1}<\frac{\check{a}}{\sqrt{\check{a}^{2}+(\hat{a}-\check{a})^{2}}} .
$$

Notice that condition (4.7) is satisfied by jumping nonlinearities:

$$
f(x, u)=\hat{a} u^{+}-\check{a} u^{-}+g(x, u)
$$

with $g$ sublinear.
Our method is useful not only for selfadjoint problems considered in Hilbert spaces: $L^{2}$ or Sobolev spaces. The boundary value problems for ordinary differential equations admit Green operators which are described sufficiently well. This enables us to divide this operator into two parts: $G_{0}$ and the other summand, and to consider the problem in different spaces. Sometimes, a perturbation of $L\left(\lambda_{0}\right)$ is not expressed by the identity operator: see [22] and [23] where the problems of the form

$$
\begin{aligned}
& \sum_{j=0}^{m} a_{j}(\lambda) x^{(j)}=f(t, x) \\
& \sum_{j=0}^{m-1}\left(b_{i j} x^{(j)}(0)+c_{i j} x^{(j)}(1)\right)=0, \quad i=1, \ldots, m
\end{aligned}
$$

are studied. The functions $a_{j}$ are analytic in a nbhd of $\lambda_{0}$ and the corresponding linear homogeneous problem has nontrivial solutions.

The perturbation method can be applied to differential equations in Banach spaces. The following problem was studied in [27]:

$$
\begin{equation*}
x^{\prime}=A(t) x+f(t, x), \quad B_{1} x(0)+B_{2} x(1)=0 \tag{4.8}
\end{equation*}
$$

where $A:[0,1] \rightarrow L(E), f:[0,1] \times E \rightarrow E$ were continuous functions, $E$ a Banach space, $f(t, \cdot)$ compact, $f(\cdot, x)$ uniformly continuous with respect to $x$ on bounded subsets, and $B_{1}, B_{2} \in L(E)$ the space of linear bounded operators on $E$. Suppose that the operator $B_{1}+B_{2} U(1)$, where $U(t), t \in$ [ 0,1 ], is the resolvent for $x^{\prime}=A(t) x$, is a Fredholm operator of index zero in $E, B_{1}+e^{\lambda} B_{2} U(1)$ is invertible for small $|\lambda|$ and the limits

$$
h_{j}=\lim _{\lambda \rightarrow 0} \frac{\left(B_{1}+e^{\lambda} B_{2} U(1)\right) x_{j}}{\left\|\left(B_{1}+e^{\lambda} B_{2} U(1)\right) x_{j}\right\|}
$$

exist for $j=1, \ldots, n$, where $x_{1}, \ldots, x_{n}$ form a basis of $\operatorname{ker}\left(B_{1}+B_{2} U(1)\right)$. Let $f$ be sublinear. If we denote by $v_{j}, j=1, \ldots, n$, the linear continuous functional on $E$ defined by

$$
\left\langle v_{j}, h_{i}\right\rangle=\delta_{i j}, \quad v_{j} \mid\left(B_{1}+B_{2} U(1)\right)(E)=0,
$$

then the L-L condition will have the form: for any $\left(x_{\nu}\right) \subset C([0,1], E)$ with $\left\|x_{\nu}\right\| \rightarrow \infty$ and $\left\|x_{\nu}\right\|^{-1} x_{\nu} \rightarrow \sum d_{j} U(\cdot) x_{j}$, there exists $j_{0}$ such that

$$
\liminf _{\nu \rightarrow \infty} d_{j_{0}}\left\langle v_{j_{0}}, B_{2} U(1) \int_{0}^{1} U^{-1}(s) f\left(s, x_{\nu}(s)\right) d s\right\rangle>0
$$

(or limsup... $<0$ ). The boundary condition can also involve nonlinear operators and both nonlinearities can have a linear growth. The suitable assumptions are much complicated (see [27]).
5. The case of the infinite-dimensional kernel. The perturbation method admits generalizations in many directions. Below, we present one of them and its applications to boundary value problems for infinitedimensional ordinary differential equations. If one examines the boundary value problems such as

$$
\begin{gathered}
x^{\prime}=f(t, x), \quad x(0)=x(1), \\
x^{\prime \prime}+m^{2} x=f(t, x), \quad x(0)=x(\pi)=0,
\end{gathered}
$$

where $x \in E$, an infinite-dimensional Banach space, then the spaces of solutions of the homogeneous linear problems are isomorphic to $E$, so we cannot apply the above theory. Although the corresponding linear operator is not Fredholm, its behaviour is similar to the finite-dimensional case.

Let $X, Y, Z$ be Banach spaces, $J: Y \rightarrow X$ an injective bounded linear operator, and $L(\lambda): Y \rightarrow Z, \lambda \in \mathbb{R}$, a continuous family of bounded linear
operators which are invertible except $L\left(\lambda_{0}\right)$. Assume that

$$
\operatorname{ker} L\left(\lambda_{0}\right)=J^{-1}\left(E_{1}\right) \oplus \ldots \oplus J^{-1}\left(E_{n}\right)
$$

where $E_{1}, \ldots, E_{n}$ are reflexive subspaces of $X$. Moreover, let

$$
J L(\lambda)^{-1}=: G(\lambda)=G_{0}(\lambda)+\sum_{j=1}^{n} c_{j}(\lambda) u_{j}(\lambda)
$$

where $G_{0}(\lambda): Z \rightarrow X$ and $u_{j}(\lambda): Z \rightarrow E_{j}, j=1, \ldots, n$, are bounded linear operators having continuous extensions to $\lambda=\lambda_{0}$, and $\left|c_{j}(\lambda)\right| \rightarrow \infty$ as $\lambda \rightarrow \lambda_{0}$. Assume that there exists a closed subspace $\widetilde{X} \subset X$ such that $G_{0}(\lambda)(Z) \subset \widetilde{X}$ and

$$
\begin{gathered}
\widetilde{X} \oplus E_{1} \oplus \ldots \oplus E_{n}=X, \\
L\left(\lambda_{0}\right)(Y)=\bigcap_{j=1}^{n} \operatorname{ker} u_{j}\left(\lambda_{0}\right), \\
L\left(\lambda_{0}\right) G_{0}\left(\lambda_{0}\right) z=z \quad \text { for } z \in L\left(\lambda_{0}\right)(Y) .
\end{gathered}
$$

We look for a solution of (4.1) or, equivalently,

$$
\begin{equation*}
x=G_{0}\left(\lambda_{0}\right) N(x)+\sum_{j} d_{j}, \quad u_{j}\left(\lambda_{0}\right) N(x)=0, \quad j=1, \ldots, n \tag{5.1}
\end{equation*}
$$

where $d_{j} \in E_{j}, j=1, \ldots, n$, are arbitrary. We suppose that the nonlinear operator $N: X \rightarrow Z$ is uniformly continuous on bounded sets and, for some sequences of finite rank bounded linear projectors $\left(P_{j, k}\right)_{k \in \mathbb{N}}$ in $E_{j}$, $j=1, \ldots, n$, such that $P_{j, k} d_{j} \rightarrow d_{j}$ as $k \rightarrow \infty$, the following condition is satisfied:
(W) for each bounded set $K \subset X$ and $\varepsilon>0$, there exists $k_{0}$ such that

$$
\left\|N\left(\widetilde{x}+\sum_{j} P_{j, k} d_{j}\right)-N\left(\widetilde{x}+\sum_{j} d_{j}\right)\right\| \leq \varepsilon
$$

for $k \geq k_{0}, \widetilde{x} \in \widetilde{X}, d_{j} \in E_{j}, j=1, \ldots, n$, such that $\widetilde{x}+\sum d_{j} \in K$.
Moreover, suppose the operators $G_{0}(\lambda) N, \lambda \neq \lambda_{0}$, are compact.
We restrict ourselves to the sublinear case. For $\lambda \neq \lambda_{0}$ the equations $x=G(\lambda) N(X)$ have solutions due to the Rothe Fixed Point Theorem. If $\lambda_{m} \rightarrow \lambda_{0}$ and $x_{m}$ is such a solution for $\lambda=\lambda_{m}$, then we prove that a bounded sequence $\left(x_{m}\right)$ contains a subsequence convergent to a solution of our equation and we find a condition which excludes the case of unbounded $\left(x_{m}\right)$.

Let $\left(x_{m}\right)$ be bounded and $x_{m}=\widetilde{x}_{m}+\sum_{j} d_{j}^{m}$ where $\widetilde{x}_{m}=G_{0}\left(\lambda_{m}\right) N\left(x_{m}\right)$ and $d_{j}^{m}=c_{j}\left(\lambda_{m}\right) u_{j}\left(\lambda_{m}\right) N\left(x_{m}\right), j=1, \ldots, n$. All sequences $\left(\widetilde{x}_{m}\right),\left(d_{j}^{m}\right)$, $j=1, \ldots, n$, are bounded, thus, by the compactness of $G_{0}\left(\lambda_{m}\right) N$ and the
reflexivity of $E_{j}$, we can assume

$$
\widetilde{x}_{m} \rightarrow \widetilde{x} \in \widetilde{X}, \quad d_{j}^{m} \rightharpoonup d_{j}, \quad j=1, \ldots, n,
$$

where $\rightharpoonup$ denotes weak convergence. Applying condition (W) and the uniform continuity of $N$, it is easy to show that

$$
N\left(\widetilde{x}_{m}+\sum_{j} d_{j}^{m}\right) \rightarrow N\left(\widetilde{x}+\sum_{j} d_{j}\right)
$$

Hence $x=\widetilde{x}+\sum d_{j}$ satisfies (5.1), by the boundedness of $\left(d_{j}^{m}\right)_{m}$ and the fact that $\left|c_{j}\left(\lambda_{m}\right)\right| \rightarrow \infty$.

Theorem 3. Under the above assumptions, if, for any sequence $\left(x_{m}\right) \subset$ $X$ such that $\left\|x_{m}\right\| \rightarrow \infty$ and $\left\|x_{m}\right\|^{-1} x_{m} \rightarrow \sum d_{j}$, there exists $j_{0}$ and $a$ functional $r: E_{j_{0}} \times E_{j_{0}} \rightarrow \mathbb{R}$ weakly continuous and homogeneous in the second variable such that $r\left(d_{j_{0}}, d_{j_{0}}\right)>0$ and

$$
\begin{equation*}
\limsup _{m \rightarrow \infty} \alpha_{j_{0}} r\left(d_{j_{0}}, u_{j_{0}}\left(\lambda_{0}\right) N\left(x_{m}\right)\right)<0 \tag{5.2}
\end{equation*}
$$

where $\alpha_{j}=\operatorname{sgn} c_{j}(\lambda)$ for $\lambda$ close to $\lambda_{0}$, then equation (4.1) has a solution.
Proof. We shall see that the sequence of fixed points $\left(x_{m}\right)$ cannot be unbounded. If it is, we can assume that $\left\|x_{m}\right\| \rightarrow \infty$. Then

$$
\frac{x_{m}}{\left\|x_{m}\right\|}=G_{0}\left(\lambda_{m}\right) \frac{N\left(x_{m}\right)}{\left\|x_{m}\right\|}+\sum_{j} \frac{c_{j}\left(\lambda_{m}\right)}{\left\|x_{m}\right\|} u_{j}\left(\lambda_{m}\right) N\left(x_{m}\right)
$$

The first summand tends to 0 since $N$ is sublinear, hence the remaining summands are bounded and contain weakly convergent subsequences

$$
\frac{c_{j}\left(\lambda_{m}\right)}{\left\|x_{m}\right\|} u_{j}\left(\lambda_{m}\right) N\left(x_{m}\right) \rightharpoonup d_{j} \in E_{j}, \quad j=1, \ldots, n
$$

For $j=j_{0}$, this contradicts (5.2).
Remarks. If the nonlinear part $N$ is bounded, then the condition (5.2) should be satisfied only for sequences $x_{m}=\widetilde{x}_{m}+\sum d_{j}^{m}$ such that $\left\|\widetilde{x}_{m}\right\| \leq M$ where $M>\left\|G_{0}\left(\lambda_{0}\right)\right\| \sup \|N(x)\|,\left\|\sum d_{j}^{m}\right\| \rightarrow \infty,\left\|\sum d_{j}^{m}\right\|^{-1} d_{j}^{m} \rightarrow d_{j} \in E_{j}$, $j=1, \ldots, n$. When, moreover, the limits of $N\left(x_{m}\right)$ exist and depend only on $\sum d_{j}$ (denote them by $\left.N\left(d_{1}, \ldots, d_{n}\right)\right)$ then the sufficient condition for the solvability of (4.1) becomes simpler: for any $\left(d_{1}, \ldots, d_{n}\right)$ on the unit sphere, there exists $j$ such that

$$
\alpha_{j} r\left(d_{j}, u_{j}\left(\lambda_{0}\right) N\left(d_{1}, \ldots, d_{n}\right)\right)<0
$$

where $r: E_{j} \times E_{j} \rightarrow \mathbb{R}$ has the properties from Theorem 3. In the previous section, we have $E_{j}=\mathbb{R}$ for any $j$, with $r$ being the multiplication of real numbers.

Now, we shall show how this theorem works for two boundary value problems mentioned at the beginning of this section. Consider

$$
\begin{equation*}
x^{\prime}=f(t, x), \quad x(0)=x(T) \tag{5.3}
\end{equation*}
$$

where $f:[0, T] \times H \rightarrow H$ is a continuous function, $H$ is a Hilbert space, $f(0, \cdot)=f(T, \cdot)$ and the following condition is satisfied: for each bounded set $K \subset H$ and $\varepsilon>0$, there exists $k_{0} \in \mathbb{N}$ such that

$$
\begin{equation*}
\left\|f(t, x)-f\left(t, P_{k} x\right)\right\| \leq \varepsilon \tag{5.4}
\end{equation*}
$$

for $k \geq k_{0}, x \in K$ and $t \in[0, T]$, where $\left(P_{k}\right)$ is a fixed sequence of finite rank orthogonal projectors in $H$ strongly convergent to the identity map.

Put $X=Z=C([0, T], H), Y=\left\{y \in C^{1}([0, T], H): y(0)=y(T)\right\}$, let $J: Y \rightarrow X$ be the inclusion map, $L(\lambda) y=y^{\prime}-\lambda y, \lambda_{0}=0$, and $N(x)(t)=$ $f(t, x(t))$. Obviously $\operatorname{ker} L(0) \cong H, \widetilde{X}=\left\{x \in X: \int_{0}^{T} x(t) d t=0\right\}$. It is easy to verify that $N$ is uniformly continuous on bounded sets and satisfies condition (W). Moreover, one can calculate the Green operator:

$$
\begin{gathered}
G_{0}(\lambda) z(t)=e^{\lambda t} \int_{0}^{t} e^{-\lambda s} z(s) d s+\frac{e^{\lambda(t+T)}}{1-e^{\lambda T}} \int_{0}^{T} e^{-\lambda s} z(s) d s+\frac{1}{\lambda T} \int_{0}^{T} z(s) d s \\
c_{1}(\lambda)=-\lambda^{-1}, \quad u_{1}(\lambda) z=T^{-1} \int_{0}^{T} z(s) d s \in H
\end{gathered}
$$

where we identify the space of constant functions with $H$. The operators $G_{0}(\lambda) N$ are compact by the General Ascoli-Arzelà Theorem. $N$ will be sublinear if we assume that

$$
\lim _{\|x\| \rightarrow \infty} \sup _{t}\|f(t, x)\| /\|x\|=0
$$

Take $r$ to be the scalar product in $H$. Then problem (5.4) has a solution if, for any sequence $\left(x_{m}\right) \subset X$ such that $\left\|x_{m}\right\| \rightarrow \infty$ and $\left\|x_{m}\right\|^{-1} x_{m} \rightarrow d \in H$ uniformly, we have

$$
\limsup _{m \rightarrow \infty} \int_{0}^{T}\left(d, f\left(t, x_{m}(t)\right)\right) d t<0
$$

or

$$
\limsup _{m \rightarrow \infty} \int_{0}^{T}\left(x_{m}(t), f\left(t, x_{m}(t)\right)\right) d t<0
$$

The dual assumptions are obvious.
We shall study the infinite-dimensional Dirichlet problem:

$$
\begin{equation*}
x_{k}^{\prime \prime}+m_{k}^{2} x_{k}=f_{k}(t, x), \quad x_{k}(0)=x_{k}(\pi)=0, \quad k \in \mathbb{N} \tag{5.5}
\end{equation*}
$$

where $x=\left(x_{k}\right) \in \mathbf{1}^{p}$ with $p>1$ and $\left(m_{k}\right) \subset \mathbb{N}$. Assume that $f_{k}:[0, \pi] \times \mathbf{l}^{p} \rightarrow$ $\mathbb{R}$ are continuous and, for all $t$ and $x$,

$$
\sum_{k=1}^{\infty}\left|f_{k}(t, x)\right|^{p}<\infty
$$

Consider the projectors $P_{k}\left(x_{1}, \ldots, x_{k}, x_{k+1}, \ldots\right)=\left(x_{1}, \ldots, x_{k}, 0, \ldots\right), k \in \mathbb{N}$, in $\mathbf{l}^{p}$. Suppose that, for any $M>0$ and $\varepsilon>0$, there exists $k_{0} \in \mathbb{N}$ such that

$$
\begin{equation*}
\sum_{n=1}^{\infty}\left|f_{n}(t, x)-f_{n}\left(t, P_{k} x\right)\right|^{p} \leq \varepsilon \tag{5.6}
\end{equation*}
$$

for $k \geq k_{0}, t \in[0, \pi]$ and $\|x\| \leq M$. Put $X=Z=C\left([0, \pi], \mathrm{l}^{p}\right), Y=$ $\left\{y \in C^{2}\left([0, \pi], \mathrm{l}^{p}\right): y(0)=y(\pi)=0\right\}, J: Y \rightarrow X$ the inclusion map, $L(\lambda)\left(y_{k}\right)=\left(y_{k}^{\prime \prime}+m_{k}^{2} y_{k}+\lambda y_{k}\right), \lambda_{0}=0$ and $N(x)(t)=\left(f_{k}(t, x(t))\right)_{k}$. Then

$$
\operatorname{ker} L(0)=\left\{\left(d_{k} \sin m_{k} \cdot\right):\left(d_{k}\right) \in \mathbf{1}^{p}\right\} \cong \mathbf{1}^{p}
$$

It is easy to find the Green operator and its decomposition into a regular part $G_{0}(\lambda)$ and an irregular term. In particular,

$$
u_{1}(0) z=\left(\int_{0}^{\pi} z_{k}(s) \sin m_{k} s d s \cdot \sin m_{k} \cdot\right)_{k \in \mathbb{N}}
$$

Condition (5.6) guarantees that $N$ is uniformly continuous on bounded sets and satisfies condition (W), and the operators $G_{0}(\lambda) N$ are compact. Assume that $f=\left(f_{k}\right)$ is sublinear. The functional $r: \mathbf{1}^{p} \times \mathbf{l}^{p} \rightarrow \mathbb{R}$ which is the multiplication of the $j$ th coordinates of both vectors has the required properties. Therefore, problem (5.5) has a solution if, for any sequence $\left(x_{m}\right) \subset X$ such that $\left\|x_{m}\right\| \rightarrow \infty$ and $\left\|x_{m}\right\|^{-1} x_{m} \rightarrow\left(d_{k} \sin m_{k} \cdot\right)$ uniformly on $[0, \pi]$, there exists $j \in \mathbb{N}$ such that

$$
\limsup _{m \rightarrow \infty} d_{j} \int_{0}^{\pi} \sin m_{j} s f_{j}\left(s, x_{m}(s)\right) d s<0
$$

or

$$
\limsup _{m \rightarrow \infty} \int_{0}^{\pi} x_{m, j}(s) f_{j}\left(s, x_{m}(s)\right) d s<0
$$

(or the dual conditions $\lim \inf \ldots>0$ ). For $p=2$, we can use a weaker condition

$$
\limsup _{m \rightarrow \infty} \int_{0}^{\pi}\left(x_{m}(s), f\left(s, x_{m}(s)\right)\right) d s<0
$$

where $(\cdot, \cdot)$ is the scalar product.

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