SOLVABILITY OF OPERATOR EQUATIONS AND PERIODIC SOLUTIONS OF SEMILINEAR HYPERBOLIC EQUATIONS

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Abstract. Let A be a closed densely defined linear map in a Banach space with infinite dimensional null space and N be a nonlinear map such that A-N is pseudo A-proper. We present a solvability theory for the equation Ax-Nx=f using only the Brouwer degree theory and the finite dimensional Morse theory. Applications to the problem of existence of weak periodic solutions of (systems of) hyperbolic equations in one and several space variables are given. Both nonresonance as well as resonance problems are considered.

I. Introduction

Let H be a separable real Hilbert space, X be a Banach space density and continuously embedded in H, $A:D(A)\subset H\to H$ be a densely defined closed linear map and N be a nonlinear map from $D(A)\cap X$ into H. We are interested in studying operator equations of the form

$$(1.1) Ax - Nx = f$$

where $f \in H$ is given.

Equations of this form appear in a variety of situations, and in particular in the theory of ordinary and partial differential equations. For example, they can describe nonlinear elliptic boundary value problems, or problems concerning periodic solutions of semilinear hyperbolic equations, or Hamiltonian systems of ordinary differential equations, etc.

Eq. (1.1) has been studied extensively by various topological as well as variational methods. When A is a Fredholm map of nonnegative index, depending on the nature of a nonlinearity, various degree theories (e.g., for the compact or condensing perturbations of the identity, coincidence degree, etc.) have been used in conjunction with the Liapunov-Schmidt technique (cf. the references). When the null space of A is infinite dimensional, Eq. (1.1) is much harder to study and

the most often used topological approach is a combination of the Leray-Schauder and coincidence degrees and the monotone operator theory (cf. Brezis-Nirenberg [Br-Ni-1-2], Mawhin [Ma-1-3] and the references in there). Essential to this approach is the existence of a compact partial inverse $A^{-1}:R(A)\to H$ of A, where the range R(A) of A is closed. However, if A contains also nonzero eigenvalues of infinity multiplicity, then A^{-1} is not compact and this approach is not suitable. This is the situation that occurs when studying the existence of periodic weak solutions of semilinear wave equations in more than one space variables. We note that, using the monotone operator theory and certain approximation procedure, Amann [Am-2] was able to obtain rather general unique solvability results for Eq. (1.1) without requiring the compactness of A^{-1} .

In [Mi-4], we initiated a new approach for the study of Eq. (1.1) (with dim ker $A=\infty$) based on a Galerkin type method. This approach requires that A-N is pseudo A-proper w.r.t. a scheme $\Gamma=\{H_n,H_n,P_n\}$ for (X,H), i.e. that the corresponding Galerkin procedure leads to a solution of (1.1). Here $\{H_n\}$ is a sequence of finite dimensional subspaces of $D(A)\cap X$, whose union is dense in both X and H, and $P_n:H\to H_n$ are the orthogonal projections. The only topological tool the method requires is the Brouwer degree theory and, as shown in [Mi-4-8], it is applicable to the situations studied by the above authors as well as to many new ones when neither A^{-1} is compact nor N is of monotone type. Moreover, in studying the existence of nontrivial solutions, we also utilize the finite-dimensional Morse theory.

In this paper we shall present a solvability theory for Eq. (1.1) using the pseudo A-proper mapping approach and give applications to the problem of existence of periodic weak solutions of (systems of) hyperbolic equations in one and more space variables. Both nonresonance as well as resonance cases are considered, i.e. when the nonlinearity N stays away from the spectrum $\sigma(A)$ of A, or interacts with it in some way.

Section II is devoted to the study of nonresonance problems for Eq. (1.1) with nonlinearities N that are asymptotically close to linear maps or are asymptotically $\{B_1, B_2\}$ -quasilinear. We prove a number of solvability results for Eq. (1.1) for each f as well as the existence of nontrivial solutions of Ax + Nx = 0 when A + N is of gradient type. Resonance problems for Eq. (1.1) are studied in Section IV, using both direct and perturbation methods. The abstract theory is then used in Sections III and V to study the existence of weak periodic solutions of various classes of semilinear hyperbolic equations without resonance as well as in resonance.

II. Semilinear equations without resonance

In this part we shall study nonresonance problems for Eq. (1.1) assuming that A - N is a pseudo A-proper map. We begin by studying such problems involving not necessarily monotone quasibounded nonlinear pertubations of both selfadjoint and nonselfadjoint linear maps with possibly infinite dimensional null space. We continue our study of Eq. (1.1) involving asymptotically $\{B_1, B_2\}$ -

quasilinear nonlinearities under the nonresonance conditions introduced by Amann [Am-2]. Finally, the existence of nontrivial solutions is studied using Morse theory.

2.1 Semilinear Equations with Quasibounded Nonlinearities. We begin by defining precisely the class of (pseudo) A-proper maps. Let X and Y be separable Banach spaces, $\{X_n\}$ and $\{Y_n\}$ be finite dimensional subspaces of X and Y with $\dim X_n = \dim Y_n$ and $\operatorname{dist}(x,X_n) \to 0$ as $n \to \infty$ for each $x \in X$. If $Q_n: Y \to Y_n$ are linear projections such that $\delta = \max \|Q_n\| < \infty$, then $\Gamma = \{X_n, Y_n, Q_n\}$ is a projection scheme for (X, Y).

Let
$$D \subset X$$
 and $T: D \to Y$. Recall [Pet]

Definition 2.1 A map $T:D\subset X\to Y$ is A-proper (resp., pseudo A-proper) w.r.t. Γ if $Q_nT:D\cap X_n\to Y_n$ is continuous for each n and, whenever $\{x_{n_k}|x_{n_k}\in D\cap X_{n_k}\}$ is bounded and $Q_{n_k}Tx_{n_k}-Q_{n_k}f\to 0$ as $k\to\infty$ for some $f\in Y$, then some subsequence $x_{n_{k(1)}}\to x$ (resp., there is an $x\in D$) and Tx=f.

The pseudo A-properness of T = A - N has been established under various assumptions on A and N in [Mi-4-7] and more details with some new examples are given in [Mi-4,8]. Throughout the paper we shall assume that $(X, ||\cdot||_0)$ is a Banach space continuously and densely embedded in a Hilbert space H.

Without the monotonicity assumption on N, we have

PROPOSITION 2.1 Let X be a reflexive Banach space compactly embedded in a Hilbert space H, $A:D(A) \subset H \to H$ be a closed and densely defined linear map and $N:H\to H$ be a nonlinear continuous map. Then $A-N:D(A)\cap X\subset X\to H$ is pseudo A-proper w.r.t. $\Gamma=\{H_n,H_n,P_n\}$ for (X,H) with $P_nAx=Ax$ on H_n .

Proof. Let $\{x_{n_k} \in X_{n_k}\}$ be bounded in X and $y_k \equiv P_{n_k}(A-N)x_{n_k} \to f$ in H. Then we may assume that $x_{n_k} \to x$ (weakly) in X, $x_{n_k} \to x$ in H, and $Nx_{n_k} \to Nx$. Hence, $Ax_{n_k} = y_k + P_{n_k}Nx_{n_k} \to f + Nx$. Since A is a closed map, $x \in D(A)$ and Ax - Nx = f. \square

Let us first study Eq. (1.1) when there is no resonance at infinity. We have

THEOREM 2.1 (cf. [5, 8]) Let $A: D(A) \subset H \to H$ be a linear selfadjoint map and $\Gamma = \{H_n, H_N, P_n\}$ be a projection scheme for (X, H) with $P_n Ax = Ax$ on H_n . Let $N: X \to H$ be a nonlinear map such that for some positive constants a, b, c and r and selfadjoint $N_{\infty}: H \to H$

- (2.1) If $(A N_{\infty})x = y$ for $y \in H$, then $x \in X$ and $||x||_0 \le c||y||$.
- $(2.2) ||Nx N_{\infty}x|| \le a||x|| + b \text{ for } ||x||_0 \ge r.$
- (2.3) $0 < a < \min\{|\mu| \mid \mu \in \sigma(A N_{\infty})\}.$

Suppose that either one of the following conditions holds:

(2.4)
$$A - N : D(A) \cap X \subset X \to H$$
 is pseudo A-proper w.r.t. Γ ,

(2.5) X is a reflexive space compactly embedded in H and N is continuous in H

Then Eq. (1.1) is solvable for each $f \in H$.

Proof. In view of Proposition 2.1, $A - N : D(A) \cap X \subset X \to H$ is pseudo A-proper w.r.t. Γ if (2.5) holds. Therefore, it remains to prove the theorem assuming (2.4).

Now, since $A - N_{\infty}$ is self-adjoint in H, we have that $\min\{|\mu| \mid \mu \in \sigma(A - N_{\infty})\} < \|(A - N_{\infty})^{-1}\|^{-1}$. Let $f \in H$ be fixed and $H(t, x) = (A - N_{\infty})x - t(N - N_{\infty})x$ on $[0, 1] \times D(A) \cap X$. Then there are $\gamma > 0$, $R \ge r$ and $n \ge n_0$ such that

$$(2.6) ||P_nH(t,x) - tP_nf|| \ge \gamma \text{for} t \in [0,1], x \in \partial B_X(0,R), n \ge n_0.$$

If not, then there would exist $t_n \in [0,1], t_n \to t_0$, and $x_{n_k} \in H_{n_k}$ such that $||x_{n_k}||_0 \to \infty$ and $P_{n_k}H(t_k, x_{n_k}) - t_kP_{n_k}f \to 0$ as $k \to \infty$. Then $z_k \equiv y_{n_k} - t_kP_{n_k}(N-N_\infty)(A-N_\infty)^{-1}y_{n_k} - t_kP_{n_k}f \to 0$, where $y_{n_k} = (A-N_\infty)x_{n_k}$. By (2.1) and (2.2) $||y_{n_k}|| \ge ||x_{n_k}||_0/c \to \infty$ and

(2.7)
$$||y_{n_k}|| \le \delta ||(N - N_{\infty})(A - N_{\infty})^{-1}y_{n_k}|| + \delta ||f|| + ||z_k|| < \delta a||(A - N_{\infty})^{-1}|| ||y_{n_k}|| + \delta (b + ||f||) + ||z_k||.$$

Dividing (2.7) by $||y_{n_k}||$ and taking the limit we get $1 \leq \delta a||(A - N_{\infty})^{-1}||$ in contradiction to (2.3). Hence, (2.6) holds and the Brouwer degree

$$\deg(P_n(A-N), B_X(0,R) \cap X_n, P_n f) = \deg(P_n(A-N_\infty), B_X(0,R) \cap X_n, 0) \neq 0$$

for each $n \ge n_0$. Hence, there exists an $x_n \in B_X(0,R) \cap X_n$, such that $P_n(A-N)x_n = P_nf$ for $n \ge n_0$, and by the pseudo A-properness of A-N, there is an $x \in D(A)$ such that Ax - Nx = f. \square

Remark 2.1 If there are real numbers $\alpha < \beta$ such that $\sigma(A) \cap (\alpha, \beta)$ consists of at most a finite number of eigenvalues, and if $\lambda_k < \lambda_{k+1}$ are some consecutive eigenvalues in (α, β) and $\lambda = (\lambda_{k+1} - \lambda_k)/2$, then (2.3) holds with $N_{\infty} = \lambda I$ if $\alpha = \gamma = (\lambda_{k+1} - \lambda_k)/2$. Indeed, the spectral gap for $A - \lambda I$ induced by the gap $(\lambda_k, \lambda_{k+1})$ is $(-\gamma, \gamma)$ and therefore $(A - \lambda I)^{-1} : H \to H$ is a bounded selfadjoint map whose spectrum lies in $(-1/\gamma, 1/\gamma)$. Hence, $||(A - \lambda I)^{-1}|| = 1/\gamma$.

Remark 2.2 If X = H and $N_{\infty} = \lambda I$ with $\lambda \notin \sigma(A)$, then (2.3) implies (2.1) since $\min\{|\mu|: |\mu \in \sigma(A - \lambda I)\} = \|(A - \lambda I)^{-1}\|^{-1}$ and $\|(A - \lambda I)x\| \ge \|(A - \lambda I)^{-1}\|^{-1}\|x\|$ for all $x \in H$. Moreover, if R(A) is the orthogonal complement of ker A and X is compactly embedded in H, then it can be shown that (2.1) with $N_{\infty} = \lambda I$ implies that dim ker A is finite.

Regarding condition (2.1), we have the following result useful in applications.

LEMMA 2.1 Let $A:D(A)\subset H\to H$ be selfadjoint with the spectum $\sigma(A)$ consisting only of eigenvalues $\{\lambda_i\mid i\in I\}$ having no accumulation points and

each λ_i have a finite multiplicity j_i . Suppose that the corresponding eigenvectors $\{e_{ij}|i\in I, 1\leq j\leq j_i\}$ form a complete basis for X and H. Suppose that there is a constant $c_0>0$ such that if Ax=y for some $y\in\ker A^\perp$, then $x\in X$ and $||x||_0\leq c_0||y||$. Then condition (2.1) holds.

Proof. For $x \in D(A)$ we have

$$x = \sum_{i} \sum_{j=1}^{j_i} x_{ij} e_{ij}$$
 and $Ax = \sum_{i} \lambda_i \sum_{j=1}^{j_i} x_{ij} e_{ij}$,

where $x_{ij} = (x, e_{ij})$. If $\lambda \notin \sigma(A)$ and $(A - \lambda I)x = y$ for $y \in H$ then

$$(A - \lambda I)x = \sum_{i} (\lambda_i - \lambda) \sum_{j=1}^{j_i} x_{ij} e_{ij} = \sum_{i} \sum_{j=1}^{j_i} y_{ij} e_{ij}$$

and therefore

$$x = (A - \lambda I)^{-1} y = \sum_{i} \frac{1}{(\lambda_i - \lambda)} \sum_{j=1}^{j_i} y_{ij} e_{ij}.$$

Since $\ker(A - \lambda I) = \{0\}$, there is a constant $\alpha > 0$ such that $|\lambda_i - \lambda| \ge \alpha$ for all $i \in I$. If not, then a subsequence $\lambda_{i_k} \to \lambda$ in contradiction to the fact that $\{\lambda_i\}$ has no accumulation points. Since $x_{ij} = y_{ij}/(\lambda_i - \lambda)$, we have that

$$|x_{ij}| \le \alpha^{-1} |y_{ij}|$$
 and $||x|| \le c_1 ||y||$

for some $c_1 > 0$ independent of y. Moreover, $x \in X$ and $||x||_0 \le c_0 ||y + \lambda x||$ if $Ax = y + \lambda x \ne 0$. If $y + \lambda x = 0$, then $x = y/\lambda \in \ker A$ and $||x||_0 = ||y||_0/\lambda \le c||y||/\lambda$ since dim ker $A < \infty$. Hence $||x||_0 \le c(\lambda)||y||$ with some $c(\lambda) > 0$. \square

If A is not selfadjoint, analyzing the proof of Theorem 2.1 we see that the following more general version of it holds.

THEOREM 2.2 Let $C: X \to H$ be a linear map such that $A-C: D(A) \cap X \to H$ is a bijection and for some positive constants $a, b, c, \delta = \max ||P_n||$ and r with $ac\delta < 1$:

(2.8)
$$||(A-C)^{-1}y||_0 \le c||y|| for each y \in H,$$

$$||Nx - Cx|| \le a||x||_0 + b \quad \text{for each} \quad ||x||_0 \ge r.$$

Then the conclusions of Theorem 2.1 are valid.

2.2 Semilinear Equations with $\{B_1, B_2\}$ -quasilinear Perturbations. We continue our study of Eq. (1.1) with the so-called asymptotically $\{B_1, B_2\}$ -quasilinear nonlinearities N. Such maps have been introduced by Perov [Per] and Krasnoselskii-Zabreiko [Kr-Za] in their study of the existence of fixed points of compact maps. Let $B_1, B_2 : H \to H$ be selfadjoint maps with $B_1 \leq B_2$, i.e. $(B_1x, x) \leq (B_2x, x)$ for $x \in H$.

Definition 2.2 a) A nonlinear map $K: H \to H$ is $\{B_1, B_2\}$ -quasilinear on a set $S \subset H$ if for each $x \in S$ there exists a linear selfadjoint map $B: H \to H$ such that $B_1 \leq B \leq B_2$ and Bx = Kx;

b) A map $N: H \to H$ is said to be asymptotically $\{B_1, B_2\}$ -quasilinear if there is a $\{B_1, B_2\}$ -quasilinear outside some ball map K such that

$$|N-K| = \limsup_{\|x\| \to \infty} \frac{\|Nx - Kx\|}{\|x\|} < \infty.$$

We do not require that $\{B_1, B_2\}$ is a regular pair as in [Kr-Za, Per]. This class of maps is rather large. For example, let $K: H \to H$ have a selfadjoint weak Gateaux derivative N' i.e. N'(x) is a selfadjoint map on H for each x and

$$\lim_{t \to 0} t^{-1} (N(x+th) - N(x), y) = (N'(x)h, y)$$

for all $x, y, h \in H$. Assume that $B_1 \leq N'(x) \leq B_2$ for each x and some selfadjoint maps B_1 and B_2 . Then N is asymptotically $\{B_1, B_2\}$ -quasilinear with |N-K| = 0. Indeed, for every $x, y, z \in H$, the mean value theorem implies the existence of a number $t \in (0, 1)$ such that

$$(Nx-Ny,z)=\left(N'(y+t(x-y))(x-y),z\right).$$

When y = 0, this gives

$$(Nx - N'(tx)x, z) = (N(0), z)$$

and therefore, if we set Kx = N'(tx)x, then

$$||Nx - Kx|| = \sup_{\|z\| \le 1} |(N(0), z)| \le ||N(0)||.$$

Hence, |N-K|=0 and N is asymptotically $\{B_1,B_2\}$ -quasilinear. In the nondifferentiable case, if Nx=B(x)x+Mx for some nonlinear map M with the quasinorm $|M|<\infty$ and selfadjoint maps $B(x):H\to H$ with $B_1\leq B(x)\leq B_2$ for each $x\in H$, then N is asymptotically $\{B_1,B_2\}$ -quasilinear.

We need the following preliminary result (cf. [Mi-8]).

LEMMA 2.2 Let $A:D(A)\subset H\to H$ and $B^\pm\in L(H)$ be selfadjoint with $B^-\leq B^+$ and H^\pm be subspaces with $H=H^-\oplus H^+$ and such that for some $\gamma_1>0$ and $\gamma_2>0$

(2.10)
$$((A - B^{-})x, x) \leq -\gamma_{1} ||x||^{2} \text{ for all } x \in H^{-} \cap D(A),$$

$$(2.11) ((A - B^+)x, x) \ge \gamma_2 ||x||^2 for all x \in H^+ \cap D(A).$$

Then there are $\epsilon > 0$ and c > 0 such that for any selfadjoint maps $B_1, B_2, C \in L(H)$ with $B_1 \leq B^-$ and $B^+ \leq B_2$ and

$$(2.12) B_1 - \epsilon I \le C \le B_2 + \epsilon I$$

we have that

$$(2.13) ||Ax - Cx|| \ge c||x|| for x \in D(A).$$

Proof. If (2.13) does not hold, then there would exist selfadjoint maps B_{1n} , B_{2n} and $C_n \in L(H)$ and $x_n \in D(A)$ with $||x_n|| = 1$ such that $B_{1n} \leq B^-$, $B^+ \leq B_{2n}$ and

$$(2.14) B_{1n} - I/n \le C_n \le B_{2n} + I/n$$

and

$$||Ax - C_n x_n|| \le ||x_n||/n.$$

Set $y_n = Ax_n - C_nx_n$. Then, we have $x_n = x_{1n} + x_{2n} \in H^- \oplus H^+$ and

$$(Ax_{2n}-C_nx_{2n},x_{2n})-(Ax_{1n}-C_nx_{1n},x_{1n})=(y_n,x_{2n}-x_{1n}).$$

Since $B_{1n} \leq B^-$ and $B^+ \leq B_{2n}$, (2.14) and (2.12)-(2.13) imply that

$$(Ax_{1n}-C_nx_{1n},x_{1n})\leq (Ax_{1n}-B^-x_{1n},x_{1n})+||x_{1n}||^2/n\leq (1/n-\gamma_1)||x_{1n}||^2.$$

Subtracting the first equation from the second one, we get

$$\gamma_{1}||x_{1n}||^{2} + \gamma_{2}||x_{2n}||^{2} - (||x_{1n}||^{2} + ||x_{2n}||^{2})/n
\leq (Ax_{2n} - C_{n}x_{2n}, x_{2n}) - (Ax_{1n} - C_{n}x_{1n}, x_{1n})
= (y_{n}, x_{2n} - x_{1n} \leq ||x_{2n} - x_{1n}||/n.$$

Hence, if $\gamma = \min\{\gamma_1, \gamma_2\}$, by the parallelogram law we get

$$(\gamma - 1/n) (||x_{2n}||^2 + ||x_{1n}||^2) \le 2/n (||x_{2n}||^2 + ||x_{1n}||^2)$$

and therefore $\gamma \leq 1/n \to 0$ as $n \to \infty$, in contradiction to $\gamma > 0$. Hence, (2.13) holds. \square

Now we are ready to prove our basic solvability result for Eq. (1.1) involving asymptotically $\{B_1, B_2\}$ -quasilinear nonlinearities N without resonance at infinity. It is based on the following continuation theorem.

THEOREM 2.3 [Mi-3] Let V be dense subspace of a Hilbert space H, $D \subset H$ be open and bounded subset and a homotopy $H: [0,1] \times (\bar{D} \cap V) \to H$ be such that

- (i) H is an A-proper homotopy w.r.t. $\Gamma = \{H_n, P_n\}$ on $[0, \epsilon] \times (\partial D \cap V)$ for each $\epsilon \in (0, 1)$ and H_1 is pseudo A-proper w.r.t. Γ ,
 - (ii) H(t,x) is continuous at 1 uniformly for $x \in \partial D \cap V$,
 - (iii) $H(t,x) \neq f$ and $H(0,x) \neq tf$ for $t \in [0,1]$, $x \in \partial D \cap V$.
 - (iv) $deg(P_nH_0, D \cap H_n, 0) \neq 0$ for all large n.

Then the equation H(1,x) = f is solvable in $\bar{D} \cap V$.

Now, we have (cf. [Mi-8])

THEOREM 2.4 Let $A:D(A)\subset H\to H$ satisfy (2.10)–(2.11) and $N:H\to H$ be bounded asymptotically $\{B_1,B_2\}$ -quasilinear with |N-K| sufficiently small. Suppose that a selfadjoint map $C_0\in L(H)$ satisfies (2.12) with sufficiently small ϵ and that $H(t,.)=A-(1-t)C_0-tN$ is A-proper w.r.t. $\Gamma=\{H_n,P_n\}$ with $P_nAx=Ax$ on H_n for each $t\in[0,1)$ and H_1 is pseudo A-proper w.r.t. Γ . Then Eq. (1.1) is solvable for each $f\in H$.

Proof. Since $N_f x = Nx - f$, $f \in H$, has the same properties as N, it suffices to solve the equation Ax - Nx = 0. Let $\epsilon_0 > 0$ be such that $|N - K| + \epsilon_0 < c$, where c is from Lemma 2.2. Then there is an r > 0 such that

$$||Nx - Kx|| \le (|N - K| + \epsilon_0)||x||$$
, for each $||x|| \ge r$.

Moreover,

$$(2.15) \quad H(t,x) = Ax - (1-t)C_0x - tNx \neq 0 \quad \text{for} \quad x \in \partial B_r \cap D(A), t \in [0,1].$$

If not, then H(t,x) = 0 for some ||x|| = r and $t \in [0,1]$. Hence, subtracting tKx from both sides, we get

$$||Ax - tKx - (1-t)C_0x|| = t||Nx - Kx|| < c||x||.$$

Since K is $\{B_1, B_2\}$ -quasilinear, there is a selfadjoint map $C^* \in L(H)$ such that $Kx = C^*x$, $B_1 \leq C^* \leq B_2$ and therefore

But, $C = tC^* + (1-t)C_0$ is selfadjoint, satisfies (2.12) in Lemma 2.1 and therefore (2.13) holds. This contradicts (2.16) and so (2.15) is valid.

Next, since C_0 and N are bounded maps, H(t,x) is an A-proper homotopy on $[0,\epsilon]\times(\bar{B}_r\cap D(A))$ w.r.t. Γ for each ϵ in (0,1) and is continuous at 1 uniformly for $x\in\bar{B}_r\cap D(A)$. Hence, the solvability of Ax-Nx=0 follows from Theorem 2.3. \square

Let us now discuss some conditions on B^{\pm} which imply (2.10)-(2.11). Assume, as in Amann [Am-2],

- (2.17) a) $A: D(A) \subset H \to H$ is self adjoint
 - b) $B^{\pm} = \sum_{i=1}^{m} \lambda_{i}^{\pm} P_{i}^{\pm}$ commute with A, where $P_{i}^{\pm} : H \to \ker(B^{\pm} \lambda_{i})$ are orthogonal projections, $\lambda_{1}^{\pm} \le \cdots \le \lambda_{m}^{\pm}$ and λ_{i}^{\pm} are pairwise distinct.
 - c) $\bigcup_{i=1}^{m} \left[\lambda_i^-, \lambda_i^+ \right] \subset \rho(A)$ -the resolvent set of A.

Being selfadjoint, A possesses a spectral resolution

$$A = \int_{-\infty}^{\infty} \lambda dE_{\lambda}$$

where $\{E_{\lambda} \mid \lambda \in R\}$ is a right continuous spectral family. Since B^{\pm} commute with A, it is known that P_i^{\pm} commute with the resolution of the identity $\{E_{\lambda} \mid \lambda \in R\}$. Hence, the selfadjoint maps $A-B^{\pm}$ have the spectral resolution

(2.18)
$$A - B^{\pm} = \sum_{i=1}^{m} \int_{-\infty}^{\infty} (\lambda - \lambda_i^{\pm}) dE_{\lambda} P_i^{\pm}.$$

Define the orthogonal projections P^{\pm} by

$$P^{-} = \sum_{i=1}^{m} E(-\infty, \lambda_{i}^{-}) P_{i}^{-}$$
 and $P^{+} = \sum_{i=1}^{m} E(\lambda_{i}^{+}, \infty) P_{i}^{+}$,

where

$$E(\alpha,\beta) = \int_{\alpha}^{\beta} dE_{\lambda}$$

for all $\alpha, \beta \in \rho(A) \cup \{\pm \infty\}$ with $\alpha < \beta$. Define $H^{\pm} = P^{\pm}(H)$ and note that by (2.17)-c,

$$P^+ = \sum_{i=1}^m E(\lambda_i^-, \infty) P_i^+,$$

and

$$\gamma = \operatorname{dist}\left(\bigcup_{i=1}^{m} [\lambda_i^-, \lambda_i^+], \sigma(A)\right) > 0.$$

Moreover, by (2.18), we have that

$$((A - B^{-})x, x) \le -\gamma ||x||^{2}$$
 for $x \in D(A) \cap H^{-}$,
 $((A - B^{-})x, x) \ge \gamma ||x||^{2}$ for $x \in D(A) \cap H^{+}$.

Hence, we have

LEMMA 2.3 If (2.17) holds and $P_i^- = P_i^+$ for $1 \le i \le m$, then there are orthogonal subspaces H^\pm such that $H = H^- \oplus H^+$ and conditions (2.10)-(2.11) hold with $\gamma_1 = \gamma_2 = \gamma > 0$.

Proof. It remains only to show that H^+ and H^- are orthogonal. Since $P_i^- = P_i^+$ for $i = 1, \ldots, m$, and $E(-\infty, \lambda_i^-) = I - E(\lambda_i^-, \infty)$, we get that $P^+ = I - P^{-}$ and therefore, $H^+ = (H^-)^\perp$. \square

When B^{\pm} are not of the form (2.17)-b, we need to assume more on the linear part A.

- (2.19) Suppose A is selfadjoint possessing a countable spectrum $\sigma(A)$ consisting of eigenvalues and whose eigenvectors form a complete orthonormal system in H.
- (2.20) There are selfadjoint maps $C_1, C_2 \in L(H)$ and two consecutive finite multiplicity eigenvalues $\lambda_k < \lambda_{k+1}$ of A such that

$$|\lambda_k||x||^2 < (C_1x, x) \le (C_2x, x) < |\lambda_{k+1}||x||^2 \quad \text{for } x \in H \setminus \{0\}.$$

Let H^- (resp. H^+) be the subspaces of H spanned by the eigenvectors of A corresponding to the eigenvalues $\lambda_i \leq \lambda_k$ (resp. $\lambda_i \geq \lambda_{k+1}$).

LEMMA 2.4 (cf. [Mi-8]) Let (2.19)–(2.20) hold. Then there are $\gamma_1 > 0$ and $\gamma_2 > 0$ such that for any selfadjoint maps $B^{\pm} \in L(H)$ satisfying $C_1 \leq B^-$ and $B^+ \leq C_2$ on H, we have that (2.10)–(2.11) hold.

Proof. It is enough to show that (2.10) holds since the same arguments also give (2.11). Since $\lambda_k < C_1 \le B^-$, it is enough to show that there is a $\gamma_1 > 0$ such that

$$(Ax - C_1x, x) \le -\gamma_1 ||x||^2 \quad \text{for all} \quad x \in D(A) \cap H^-.$$

If such a $\gamma_1 > 0$ did not exist, then there would exist $\{x_n\} \subset D(A) \cap H^-$ with $||x_n|| = 1$ and such that

$$-1/n \leq (Ax_n - C_1x_n, x_n) \leq (Ax_n - \lambda_k x_n, x_n), \quad n = 1, 2, \dots$$

Decompose $H^- = \tilde{H} \oplus \bar{H}$, where \tilde{H} is spanned by the eigenvaectors $\{e_i\}$ corresponding to $\lambda_i \leq \lambda_{k-1}$ and \bar{H} is the finite dimensional space spanned by the eigenvectors corresponding to λ_k . Then $x_n = \tilde{x}_n + \bar{x}_n \in \tilde{H} \oplus \bar{H}$ and

$$-1/n \le (Ax_n - \lambda_k x_n, x_n) = (A\tilde{x}_n - \lambda_k \tilde{x}_n, \tilde{x}_n)$$

$$= \sum_{i \le k-1} \lambda_i (\tilde{x}_n, e_i)^2 - \lambda_k ||\tilde{x}_n||^2 \le (\lambda_{k-1} - \lambda_k) ||\tilde{x}_n||^2$$

Hence, $\tilde{x}_n \to 0$ as $n \to \infty$ and $x_n \to \bar{x} \in \bar{H}$ with $||\bar{x}|| = 1$ since $||x_n||^2 = 1 - ||\tilde{x}_n||$ and \bar{H} is finite dimensional. Thus,

$$-1/n \le (\lambda_k x_n - C_1 x_n, x_n) + (A x_n - \lambda_k x_n, x_n) \le (\lambda_k x_n - C_1 x_n, x_n) + (\lambda_{k-1} - \lambda_k) ||\tilde{x}_n||^2$$

and passing to the limit as $n \to \infty$, we get

$$0 \leq ((\lambda_k - C_1)\bar{x}, \bar{x}), \bar{x} \neq 0$$

in contradiction to (2.20). Hence, (2.10) holds. \square

Remark 2.3 If λ_k (resp., λ_{k+1}) is of infinite multiplicty, then Lemma 2.3 is still valid if we assume in (2.20)

$$(\lambda_k + \epsilon)||x||^2 \le (C_1 x, x)$$
 resp. $(C_2 x, x) \le (\lambda_{k+1} - \epsilon ||x||^2)$ for $0 \ne x \in H$.

This is easy to check by analyzing the proof of this lemma.

Next, let us look at the case when $H^- \oplus H^+ \neq H$, which is also useful in applications.

Recall that a closed subspace $X \subset H$ is said to reduce A if A commutes with the orthogonal projection P of H onto X, i.e., if PA = AP. As before, let P^{\pm} be the orthogonal projections of H onto H^{\pm} . Regarding A and H we require that the following conditions hold for a scheme $\Gamma = \{H_n, P_n\}$ (cf. [Am-2]):

- (2.21) (i) $H_1 \subset H_2 \subset \cdots$ with each H_n closed in H, dim $H_n = \infty$ and $\bigcup H_n$ is dense in H;
 - (ii) each H_n reduces A;
 - (iii) the orthogonal projections $P_n: H \to H_n$ commute with P^{\pm} ;
 - (iv) $H_n = H_n^- \oplus H_n^+$, where $H_n^{\pm} = H^{\pm} \cap H_n$;

(v) $Q^{\pm}(D(A) \cap H_n) \subset D(A)$ for each n, where $Q^{\pm}: H^- \oplus H^+ \to H^{\pm}$ are orthogonal projections.

We have the following extension of Theorem 2.4.

Theorem 2.5 [Mi-8] Let H^{\pm} be closed subspaces of H such that $H^{-} \cap H^{+} = \{0\}$, $A: D(A) \subset H \to H$ satisfy conditions (2.10)–(2.11) and (2.21) hold. Suppose that $N: H \to H$ is bounded asymptotically $\{B_1, B_2\}$ -quasilinear with |N-K| sufficiently small and $A-N: D(A) \subset H \to H$ is pseudo A-proper w.r.t. $\Gamma = \{H_n, P_n\}$. Suppose that a selfadjoint map $C_0 \in L(H)$ satisfies (2.12) and for each large n the map $H_n(t, \cdot) = P_n(A - (1-t)C_0 - tN): D(A) \cap H_n \to H_n$ is A-proper w.r.t. $\Gamma_n = \{H_{n,k}, Q_k\}$ for H_n , with $Q_kAx = Ax$ on $H_{n,k}$, for each $t \in [0,1)$ and $H_n(1, \cdot)$ be pseudo A-proper w.r.t. Γ_n . Then Eq. (1.1) is solvable for each $f \in H$.

When $P_i^- \neq P_i^+$ for $1 \leq i \leq m$ and $H \neq H^- \oplus H^+$, regarding conditions (2.10)–(2.11) we have (cf. [Am-2]).

LEMMA 2.5 Let (2.17) hold and there exist a unitary map $U \in L(H)$ such that A commutes with U and

$$P_i^- = UP_i^+U^{-1}, \quad i = 1, \dots, m.$$

Suppose that A has a pure point spectrum in $(\lambda_1^-, \lambda_m^+)$. Then there are closed subspaces H^{\pm} of H with $H^{-} \cap H^+ = \{0\}$ such that conditions (2.10)-(2.11) hold with $\gamma_1 = \gamma_2 = \gamma > 0$, and a scheme $\Gamma = \{H_n, P_n\}$ satisfying (2.21).

2.3 Asymptotically linear equations and Morse theory. In this section we shall study the existence of nontrivial solutions in a Hilbert space H of equations of the form

$$(2.22) Ax + Nx = 0$$

of potential type using a combination of the finite dimensional Morse theory and the A-proper mapping approach. We note that Galerkin methods in conjunction with Morse theory have been used earlier by many authors in the study of equations of the form (2.22) with N compact.

The novelty in our approach is that we require that A + N is A-proper with respect to a projectionally complete scheme, which is so in particular when N is compact, or ball condensing or monotone (of Section 2 in [Mi-8] and Proposition 2.3).

Throughout the section, we assume that $A:D(A)\in H\to H$ is a linear densely defined, selfadjoint map with dim ker $A\leq \infty$ and $N:H\to H$ is a nonlinear map such that $A+N:D(A)\subset H\to H$ is potential and A-proper with respect to a projectinally complete scheme $\Gamma=\{H_n,P_n\}$ with $P_nAx=Ax$ for $x\in H_n$.

Let $f: H \to R$ be such that f'(x) = Ax + Nx. We say that f is asymptotically quadratic at infinity if there is a selfadjoint bounded linear map $N_{\infty}: H \to H$ such that

$$(2.23) |N - N_{\infty}| = \limsup_{\|x\| \to \infty} \frac{\|Nx - N_{\infty}x\|}{\|x\|} < \infty.$$

Definition 2.3 A C^1 function $f: H \to R$ is said to satisfy the Palais-Smale condition (PS) if any sequence $\{x_n\} \in H$ with $f'(x_n) \to 0$ and $\{f(x_n)\}$ bounded possesses a convergent subsequence.

Let f_n be the restriction of f to H_n .

Definition 2.4 f is said to satisfy the $(PS)_n$ condition if any sequence $\{x_n|x_n \in H_n\}$ with $||f'(x_n)|| \to 0$ and $\{f(x_n)\}$ bounded possesses a convergent subsequence.

The next result provides conditions on the A-proper map f' = A + N which guarantee condition $(PS)_n$ for f.

PROPOSITION 2.2 Let $f: H \to R$ with $f_n \in C^1(H_n, R)$, (2.23) hold with $|N - N_{\infty}|$ sufficiently small, $A + N_{\infty}$ and A + N be A-proper with respect to $\Gamma = \{H_n, P_n\}$ and $A + N_{\infty}$ be injective. Then

(i) f_n is asymptotically quadratic and

$$\limsup_{\|x\|\to\infty, x\in H_n} \frac{\|f'(x)-(A+P_nN_\infty)x\|}{\|x\|} \leq |N-N_\infty|.$$

- (ii) f_n satisfies condition (PS) for each large n.
- (iii) f satisfies condition (PS)_n.
- (iv) The equation Ax + Nx = f is approximation solvable for each $f \in H$.

Proof. (i) Let $\epsilon>0$ and R>0 such that $||Nx-N_{\infty}x||\leq (|N-N_{\infty}|+\epsilon)||x||$ for all $||x||\geq R$. Since $f_n'(x)=Ax+P_nNx$ for each $x\in H_n$, we get that

$$(2.24) ||f_n'(x) - (A + P_n N_\infty)x|| = ||P_n(N - N_\infty)x|| \le (|N - N_\infty| + \epsilon)||x||$$

for all $x \in H_n$, $||x|| \ge R$ which proves (i).

(ii) Since $A + N_{\infty}$ is injective and A-proper, arguing by contradiction we see that there are c > 0 and $n_0 \ge 1$ such that

(2.25)
$$||P_n(A+N_\infty)x|| \ge c||x||$$
 for all $x \in H_n$, $n \ge n_0$.

Let $\{x_k\} \subset H_n$ with $||f_n'(x_k)|| \to 0$ as $k \to \infty$ and $\{f_n(x_k)|k \ge 1\}$ bounded. Then, by (2.24) and (2.25),

$$c||x_k|| \le ||P_n(A + N_\infty)x_k|| = ||P_n(A + N)x_k - P_n(N - N_\infty)x_k||$$

$$= ||f'_n(x_k) - P_n(N - N_\infty)x_k|| \le ||f'_n(x_k)|| + ||f'_n(x_k) - (A + P_nN_\infty)x_k||$$

$$\le \epsilon + (|N - N_\infty| + \epsilon)||x_k||, \quad \text{for all} \quad n \ge n_1 \ge n_0.$$

Taking $\epsilon > 0$ such that $|N - N_{\infty}| + \epsilon < c$, this yields the boundedness of $\{x_k\}$ in H_n . Hence a subsequence $\{x_{k_i}\}$ converges in H_n .

(iii) Let $\{x_n \in H_n\}$ satisfy $||f'_n(x_n)|| \to 0$. As above, by (2.24) and (2.25), the sequence $\{x_n\}$ is bounded. Since $f'_n(x_n) = Ax_n + P_n Nx_n \to 0$, by the A-properness

of A + N a subsequence $x_{n_k} \to x$ and f'(x) = Ax + Nx = 0, i.e. x is a critical point of f.

- (iv) This is a special case of Theorem 2.1. \square
- Remark 2.4 The boundedness of $\{f_n(x_k) | k \ge 1\}$ and $\{f(x_n)\}$ was not used in the proofs of (ii) and (iii).

Next, we give special classes of A-proper maps appearing in Proposition 2.2.

PROPOSITION 2.3 Let $A:D(A)\subset H\to H$ be a densely defined Fredholm map of index zero and $N:H\to H$ be a continuous nonlinear and bounded map. Suppose that either one of the following conditions holds:

- (i) N is k-ball contractive with k sufficiently small and $|N N_{\infty}| = 0$ for some continuous linear map N_{∞} ;
 - (ii) The partial inverse A^{-1} of A is compact.

Then $A + N : D(A) \subset H \to H$ and $A + N_{\infty}$ are A-proper with respect to $\Gamma = \{H_n, P_n\}$ with $P_n Ax = Ax$ for $x \in H_n$.

Proof (i) Since the index i(A) = 0, we have that $H = H_0 \oplus H_1$, where $H_0 = \ker A$. Let $\{x_{n_k} \in H_{n_k}\}$ be bounded and $y_k \equiv P_{n_k}(A+N)x_{n_k} \to f$. Then $x_{n_k} = x_{0n_k} + x_{1n_k}$ with $x_{in_k} \in H_i$, i = 0, 1, and we may assume that $x_{0n_k} \to x_0 \in H_0$ and

$$x_{1n_k} + A^{-1}(I-P)P_{n_k}Nx_{n_k} = A^{-1}(I-P)y_k \to A^{-1}(I-P)f,$$

where $P: H \to H_0$ is the orthogonal projection. Hence, the ball-measure of noncompactness

$$\chi(\{x_{1n_k}\}) = \chi(\{A^{-1}(I-P)P_{n_k}Nx_{n_k}\}) \le ||A^{-1}||k\chi(\{x_{n_k}\}) = k||A^{-1}||\chi(\{x_{1n_k}\}).$$

Since k is sufficiently small, we get that $\chi(\{x_{1n_k}\})=0$ and therefore $x_{1n_k}\to x_1\in H_1$. Thus, $x_{n_k}\to x=x_0+x_1$ and Ax+Nx=f. Since it is well known that N_∞ is also k-ball contractive, the same arguments show that $A+N_\infty$ is also A-proper with respect to Γ .

(ii) Let $\{x_{n_k}\}$ be as in i). Then $x_{0n_k} \to x_0 \in H_0$ and $P_{n_k}Nx_{n_k} \to u$ by the boundedness of N. Hence, $x_{1n_k} = A^{-1}(I-P)(y_k-P_{n_k}Nx_{n_k}) \to A^{-1}(I-P)(f-u) = x_1$ so that $x_{n_k} \to x = x_0 + x_1$. Since $Nx_{n_k} \to Nx$, it follows that $Ax_{n_k} = y_k - P_{n_k}Nx_{n_k} \to f - Nx$ and so Ax + Nx = f by the clossedness of A. \square

Remark 2.5 When dim ker $A = \infty$, various conditions on A and N that imply the A-properness of A + N have been discussed in [Mi-4-8].

To obtain nontrivial solutions of Eq. (2.22), we need to impose some additional conditions on A and N. To that end we shall use some finite dimensional results on the existence of nontrivial critical points. For a symmetric matrix S, denote by $m^-(S)$, $m^0(S)$ and $m^+(S)$ the negative, zero and positive Morse index of S. We have (cf. [A-Z, Ch, L-L])

THEOREM 2.6 Let $f: \mathbb{R}^n \to \mathbb{R}$ be a C^1 -function and A_0 and A_{∞} be symmetric matrices such that

$$\frac{\|f'(x) - A_{\infty}x\|}{\|x\|} \to 0 \quad as \quad \|x\| \to \infty$$

and

(2.27)
$$\frac{||f'(x) - A_0 x||}{||x||} \to 0 \quad as ||x|| \to 0.$$

If $m^0(A_0) = m^0(A_\infty) = 0$ and $m^-(A_0) \neq m^-(A_\infty)$, then f has at least one nontrivial critical point. Moreover, if for some r > 0

$$(2.28) ||f'(x) - A_0x|| < ||x||/(2||A^{-1}||) for ||x|| < r.$$

then the nontrivial critical points of f lie outside the ball B(0,r).

Suppose that there is a selfadjoint map $N_0: H \to H$ such that

(2.29)
$$||Nx - N_0x||/||x|| \to 0 \text{ as } ||x|| \to 0.$$

For each n, define the numbers $I^{-}(f_n, \infty) = m^{-}(P_n(A + N_{\infty})P_n)$, $I^{-}(f_n, 0) = m^{-}(P_n(A + N_0)P_n)$ and $I^{0}(f_n, 0) = m^{0}(P_n(A + N_0)P_n)$.

The Morse indices of $A + N_{\infty}$ and $A + N_0$ may be infinite and the Morse indices of the approximate operators $I^{-}(f_n, \infty)$ and $I^{-}(f_n, 0)$ converge to infinity as $n \to \infty$. We have the following extension of Theorem 2.6 to the infinite dimensional case.

THEOREM 2.7 Let $f: H \to R$ with $f_n \in C^1(H_n, R)$, (2.23) and (2.29) hold with $|N - N_{\infty}| = 0$ and $A + N_0$ and $A + N_{\infty}$ be both injective A-proper maps with respect to $\Gamma = \{H_n, P_n\}$. Suppose that $I^-(f_n, 0) \neq I^-(f_n, \infty)$ for all large n. Then f has at least one nontrivial critical point.

Proof. Since $A + N_{\infty}$ and $A + N_0$ are injective A-proper maps, there are c > 0 and $n_0 \ge 1$ such that for each $n \ge n_0$

$$(2.30) ||(A + P_n N_0)x|| \ge c||x|| \text{ and } ||(A + P_n N_\infty)x|| \ge c||x|| \text{ for all } x \in H_n.$$

Let r > 0 be so small that

$$||f'(x) - (A + N_0)x|| \le \frac{c}{2}||x||$$
 for $||x|| < r$.

Then, for each $n \ge n_0$ and $x \in H_n$ with ||x|| < r,

$$||f'_n(x) - (A_n + P_n N_0 P_n)x|| \le ||f'(x) - (A + N_0)x||$$

$$\le \frac{c}{2}||x|| \le \frac{1}{2}||(A_n + P_n N_0 P_n)^{-1}||^{-1}||x||$$

by (2.30), where A_n is the restriction of A to H_n . By Theorem (2.6), f_n has a nontrivial critical point $x_n \in H_n$ with $||x_n|| \ge r$ for all $n \ge n_0$. Since f satisfies

the $(PS)_n$ condition by Proposition 2.2, a subsequence $x_{n_k} \to x$ and x is a critical point of f since f' = A + N is A-proper, with $||x|| \ge r$. \square

Let us now show that under the conditions of Theorem 2.7, the conditions $I^-(f_n,0) \neq I^-(f_n,\infty)$ for all large n are implied by the condition $m^-(P_n(A+N_0)P_n) - m^-(A_n) \neq m^-(P_n(A+N_\infty)P_n) - m^-(A_n)$ for some sufficiently large n.

The following algebraic lemma is needed (see [L-L]).

LEMMA 2.6 Let S and T be two symmetric $n \times n$ matrices. Then $m^-(S) = m^-(S+T)$ provided that for any d>0 there is an r=r(d)>0 such that either one of the following conditions holds:

- (i) S is invertible, $||S^{-1}|| < d$ and ||T|| < r;
- (ii) $||S_1^{-1}|| < d$, ||T|| < r and TP = 0, where $P : \mathbb{R}^n \to \ker(S)$ is the orthogonal projection onto $\ker(S)$ and S_1 is the restriction of S to $(I P)\mathbb{R}^n$. The condition TP = 0 means that $\ker(S) \subset \ker(S + T)$.

LEMMA 2.7 Let $A: D(A) \subset H \to H$ be a Fredholm map of index zero and $B: H \to H$ be a continuous linear map such that A+B is A-proper with respect to $\Gamma = \{H_n, P_n\}$ with $P_nAx = Ax$ on H_n . If either

- (i) A + B is one-to-one, or
- (ii) $P_n P_0 = P_0$, where $P_0 : H \to \ker(A + B)$ is the orthogonal projection onto $\ker(A + B)$,

then $m^-(A_n + P_nBP_n) - m^-(A_n)$ is a constant for all large n.

Proof. (cf. also [L-L]) Let $F = \operatorname{span}\{H_n, H_m\}, Q: H \to F$ be the orthogonal projection and $Q_n: H \to F_n$ be the orthogonal projection onto F_n , the complement of H_n in F. Then $Q = P_n + Q_n$ and Q and Q_n commute with A since so do P_n and P_m . Assume that $F_n \neq \{0\}$ for each n. We represent the map A + B in F by the matrix W = S + T:

$$W = \begin{pmatrix} P_n(A+B)P_n & P_n(A+B)Q_n \\ Q_n(A+B)P_n & Q_n(A+B)Q_n \end{pmatrix},$$

$$S = \begin{pmatrix} A_n + P_nBP_n & 0 \\ 0 & Q_n(A+P)Q_n \end{pmatrix}, \quad T = \begin{pmatrix} 0 & P_nBQ_n \\ Q_nBP_n & Q_n(B-P)Q_n \end{pmatrix},$$

where $P: H \to \ker(A)$ is the orthogonal projection, and we used the fact that $Q_n A = AQ_n$ and $P_n Q_n = Q_n P_n = 0$. We have that P is compact, commutes with P_n and Q_n , and since A + P is invertible there are $c_1 > 0$ and $n_0 \ge 1$ such that

(2.31)
$$||(Q_n(A+P)Q_n)^{-1}|| < c_1 \text{ for all } n \ge n_0.$$

Assume first that A + B is invertible. Since it is A-proper with respect to Γ , there are c_2 and $n_2 \ge 1$ such that

$$||P_n(A+B)P_nx|| > c_2||P_nx||$$
 for all $x \in H$, $n \ge n_2$

and therefore, in H_n

(2.32)
$$||(P_n(A+B)P_n)^{-1}|| < 1/c_2 \text{ for all } n \ge n_2.$$

Since $||Q_nB|| = ||BQ_n|| \to 0$ and $||Q_nP|| \to 0$ as $n \to \infty$, we get by Lemma 2.6 that for each n large

$$(2.33) m^{-}(W) = m^{-}(S) = m^{-}(A_n + P_n B P_n) + m^{-}(Q_n (A + P) Q_n).$$

Using Lemma 2.6-(ii), we can also prove that (2.33) holds when $P_n P_0 = P_0$. In particular, taking B = 0, (2.33) becomes

$$(2.34) m^{-}(QAQ) = m^{-}(A_n) + m^{-}(Q_n(A+P)Q_n).$$

Combining (2.33) and (2.34), we get

$$m^{-}(A_n + P_n B P_n) - m^{-}(A_n) = m^{-}(W) - m^{-}(QAQ)$$

= $m^{-}(A_m + P_m B P_m) - m^{-}(A_m)$. \square

Using Lemma 2.7, we have the following special case of Theorem 2.7.

COROLLARY 2.1 Let $f: H \to R$ with $f_n \in C^1(H_n, R)$, A be Fredholm of index zero, (2.23) and (2.29) hold with $|N - N_{\infty}| = 0$ and $A + N_0$ and $A + N_{\infty}$ be both injective and A-proper with respect to $\Gamma = \{H_n, P_n\}$. Then, $m^-(A_n + P_nN_0P_n) - m^-(A_n)$ and $m^-(A_n + P_nN_{\infty}P_n) - m^-(A_n)$ are both constant for each large n. If these two constants are different, then f has at least one nontrivial critical point.

Proof. By Lemma 2.7 both $m^-(A_n + P_nN_0P_n) - m^-(H_n)$ and $m^-(A_n + P_nN_\infty P_n) - m^-(A_n)$ are constant for all large n. Moreover, if they are different constants, then $I^-(f_n, 0) \neq I^-(f_n, \infty)$ for all large n, and the conclusion follows from Theorem 2.7. \square

Remark 2.6 When $f \in C^1(H, R)$ with A and N compact and $A = A^*$, Lemma 2.7 and Corollary 2.1 were proven by Li-Liu [L-L] using different type of arguments.

Next, we shall treat the case when 0 is a degenerate critical point of f. We need the following finite-dimensional result.

THEOREM 2.8 [L-L] Let $f \in C^1(\mathbb{R}^n, \mathbb{R}) \cap C^2(B(0, d), \mathbb{R})$ and satisfy (2.26) and (2.27). Suppose that $m^0(A_{\infty}) = 0$ and $m^-(A_{\infty}) \notin [m^-(A_0), m^-(A_0) + m^0(A_0)]$. Then f has at least one nontrivial critical point. Moreover, if

$$||f'(x) - A_0|| \le \frac{1}{2} ||A_0^{-1}||^{-1} \quad \text{for } ||x|| < d,$$

where A_0^{-1} is the inverse of A_0 restricted to the range of A_0 , then at least one critical point of f lies outside the ball B(0, d/2).

We have the following extension to the infinite dimensional case.

THEOREM 2.9 Let $f: H \to R$ with $f_n \in C^1(H_n, R) \cap C^2(H_n \cap B(0, d), R)$ and satisfy (2.23) and (2.29) with $|N - N_{\infty}| = 0$. Suppose that $A + N_{\infty}$ is injective

and A-proper with respect to $\Gamma = \{H_n, P_n\}$ with $P_n Ax = Ax$ on H_n , and $A + N_0$ is A-proper with respect to Γ with $\ker(A + N_0) \subset H_n$ for all large n. If

(2.36)
$$I^{-}(f_n, \infty) \notin [I^{-}(f_n, 0), I^{-}(f_n, 0) + I^{0}(f_n, 0)]$$

for all large n, then f has at least one nontrivial critical point.

Proof. Since $H_0 = \ker(A + N_0) \subset H_n$, then $H = H_0 \oplus \tilde{H}$ for some closed subspace \tilde{H} and $H_n = H_0 \oplus H_{1n}$ with $H_{1n} = H_n \cap \tilde{H}$. Since the restriction of $A + N_0$ to \tilde{H} is injective and $A + N_0$ is A-proper, it is easy to see that there are c > 0 and $n_0 \ge 1$ such that

$$||P_n(A+N_0)P_nx|| > c||x||$$
 for $x \in H_{1n}, n \ge n_0$.

Hence, the inverses of the restrictions of $P_n(A+N_0)P_n$ to H_{1n} satisfy

$$||(P_n(A+N_0)P_n)^{-1}|| < 1/c \text{ for } n \ge n_0.$$

Let r < d/2 be sufficiently small so that

$$||f''(x) - (A + N_0)|| < c/2$$
 for $||x|| \le 2r$,

and therefore,

$$||f_n''(x) - (A_n + P_n N_0 P_n)|| \le ||f''(x) - (A + N_0)||$$

$$\le c/2 < ||(P_n (A + N_0) P_n)^{-1}||/2 \text{ for } ||x|| \le 2r, n \ge n_0.$$

Hence, from (2.36), (2.37) and Theorem 2.8, we get a nontrivial critical point $x_n \in H_n$ of f_n with $||x_n|| \ge r$ for each $n \ge n_0$. Since f satisfies the $(PS)_n$ condition, a subsequence $x_{n_k} \to x$ with $||x|| \ge r$ and f'(x) = 0 by the A-properness of A + N. \square

COROLLARY 2.2 Let $f: H \to R$ with $f_n \in C^1(H_n, R) \cap C^2(H_n \cap B(0, d), R)$, A be Fredholm of index zero and conditions (2.23) and (2.29) hold with $|N-N_{\infty}|=0$. Suppose that $A+N_{\infty}$ is injective and A-proper with respect to Γ and $A+N_0$ is A-proper with respect to Γ with $\ker(A+N_0) \subset H_n$ for all large n. Then, $m^-(A_n+P_nN_0P_n)-m^-(A_n)$ and $m^-(A_n+P_nN_\infty P_n)-m^-(A_n)$ are both constant for each large n. Moreover, if

$$(2.38) m^{-}(A_{n} + P_{n}N_{\infty}P_{n}) - m^{-}(A_{n})$$

$$\notin [m^{-}(A_{n} + P_{n}N_{0}P_{n}) - m^{-}(A_{n}),$$

$$m^{-}(A_{n} + P_{n}N_{0}P_{n}) - m^{-}(A_{n}) + m^{0}(A_{n} + P_{n}N_{0}P_{n})]$$

with n large, then f has at least one nontrivial critical point.

Proof. By Lemma 2.7(ii), $m^-(A_n + P_nBP_n) - m^-(A_n)$ is constant for all large n with $B = N_0$ and $B = N_{\infty}$. Moreover, (2.38) implies that

$$I^{-}(f_n, \infty) \notin [I^{-}(f_n, 0), I^{-}(f_n, 0) + I^{0}(f_n, 0)]$$

for all large n, and the conclusion follows from Theorem 2.9. \square

Finally, we shall consider the existence of nontrivial critical points of f when 0 is a degenerate critical point at which f has a local linking as defined below. Again we begin with the finite dimensional case first.

Definition 2.5 Let $f: X = \mathbb{R}^n \to \mathbb{R}$ be of class C^1 . We say that f has a local linking at 0 if there are subspaces Y and Z of X with X = Y + Z and positive constants a and r such that

(2.39)
$$f(y) \ge a \qquad \text{for } y \in \partial(B_r \cap Y)$$

$$f(y) \ge 0 \qquad \text{for } y \in B_r \cap Y$$

$$f(z) \le -a \qquad \text{for } z \in \partial(B_r \cap Z)$$

$$f(z) \le 0 \qquad \text{for } z \in B_r \cap Z.$$

THEOREM 2.10 ([L-L]) Let $f \in C^1(\mathbb{R}^n, \mathbb{R})$ satisfy (2.26), (2.27), (2.39) and $m^-(A_\infty) = 0$, $m^-(A_\infty) \neq \dim \mathbb{Z}$. Then f has at least one nontrivial critical point x with $||x|| \geq r$ and $|f(x)| \geq a$.

The following result involves the local linking condition in the infinite dimensional case.

THEOREM 2.11 Let $f: H \to R$ with $f_n \in C^1(H_n, R)$ and satisfy (2.23) and (2.29) with $|N - N_{\infty}| = 0$. Suppose that $A + N_{\infty}$ is injective and A-proper with respect to $\Gamma = \{H_n, P_n\}$ with $P_n Ax = Ax$ on H_n and for each large n there is a decomposition $H_n = Y_n + Z_n$ and positive constants a and r such that

(2.40)
$$f(y) \ge a \qquad \text{for } y \in \partial(B_r \cap Y_n) \\ f(y) \ge 0 \qquad \text{for } y \in B_r \cap Y_n \\ f(z) \le -a \qquad \text{for } z \in \partial(B_r \cap Z_n) \\ f(z) \le 0 \qquad \text{for } z \in B_r \cap Z_n.$$

If $I^-(f_n, \infty) \neq \dim Z_n - m^-(A_n)$ for infinitely many n, then f has at least one nontrivial critical point.

Proof. By Proposition 2.2 and (2.40), Theorem 2.10 yields critical points $x_n \in H_n$ of f_n with $||x_n|| \ge r$ and $|f_n(x_n)| \ge a$ for each large n. Since f satisfies condition $(PS)_n$ and A + N is A-proper, a subsequence $x_{n_k} \to x$ with f'(x) = Ax + Nx = 0 and $x \ne 0$. \square

COROLLARY 2.3 Let $f: H \to R$ with $f_n \in C^1(H_n, R)$ and satisfy (2.23) and (2.29) with $|N - N_{\infty}| = 0$. Suppose that $A + N_{\infty}$ is injective and A-proper with respect to $\Gamma = \{H_n, P_n\}$ with $P_n Ax = Ax$ on H_n and that there is a decomposition H = Y + Z such that $P_n Y \subset Y$ and $P_n Z \subset Z$ and for some positive constants a and r such that

$$f(y) \ge a \qquad for \ y \in \partial(B_r \cap Y)$$

$$f(y) \ge 0 \qquad for \ y \in B_r \cap Y$$

$$f(z) \le -a \qquad for \ z \in \partial(B_r \cap Z)$$

$$f(z) \le 0 \qquad for \ z \in B_r \cap Z.$$

If $I^-(f_n, \infty) \neq \dim P_n Z - m^-(A_n)$ for infinitely many n, then f has at least one nontrivial critical point.

Remark 2.7 When $f \in C^1(H,R) \cap C^2(B(0,r),R)$, respectively $f \in C^1(H,R)$, and A and N are compact maps with $A = A^*$, Corollary 2.2, respectively Theorem 2.11 and Corollary 2.3, were proven by Li-Liu [L-L] using different type of arguments.

III. Periodic Solution of Semilinear Wave Equations Without Resonance

Let $Q \subset \mathbb{R}^n$ be a bounded domain with smooth boundary and set $\Omega = (0,T) \times Q$ and let $H = L_2(\Omega,\mathbb{R}^m)$ with the inner product defined by

$$(u,v) = \int_0^T \int_Q \left(u(t,x),v(t,x)\right) dx dt$$

where (u(t,x),v(t,x)), $(t,x) \in \Omega$, is the inner product in R^m . Let L_1 be a linear selfadjoint elliptic operator in space variables $x \in R^n$ with coefficients independent of t such that the induced bilinear form a(u,v) on the Sobolev space $W_2^1(Q,R^m)$ is continuous and symmetric. Suppose that V is a closed subspace of $W_2^1(Q,R^m)$, containing the test functions, such that a(u,v) is semi-coercive on V, i.e. there are constants $a_1 > 0$ and $a_2 \ge 0$ such that

$$a(u, v) \ge a_1 ||u||_{2,1}^2 - a_2 ||u||_{L_2}^2$$
 for all $u \in V$.

Define a linear map $L_0: D(L_0) \subset L_2(Q, \mathbb{R}^m) \to L_2(Q, \mathbb{R}^m)$ by

$$(L_0u,v)=a(u,v)$$
 for each $v\in V$,

where

$$D(L_0) = \{u \in V \mid a(u, \cdot) \text{ is continuous on } V \text{ in the } L_2\text{-norm}\}.$$

It is well known that L_0 is selfadjoint and has a compact resolvent, since W_2^1 is compactly embedded in L_2 . Next, define a selfadjoint map with compact resolvent $L: D(L) \subset H_1 = L_2(Q, R^m) \to H_1$ by:

$$D(L) = [D(L_0)]^m \quad \text{and} \quad L = \operatorname{diag}(L_0, \dots, L_0).$$

Since L has a compact resolvent, there is an orthonormal basis $\{\psi_j \mid j \in J\}$ in H_1 and a sequence of its eigenvalues $\{\mu_j \mid j \in T\}$ such that $|\mu_j| \to \infty$.

Let $F: R \times Q \times R^m \to R^m$ be a Caratheodory function and consider the semilinear system of wave equations:

(3.1)
$$\begin{cases} u_{tt} - L_1 u - F(t, x, u) = f(t, x) \\ u(t, .) \in V^m \end{cases}$$

where $f \in L_2(\Omega, \mathbb{R}^m)$ is a T-periodic function in t variable and $\tau = 2\pi/T$ is rational.

By a T-periodic weak solution of the variational boundary value problem (3.1) for the semilinear system of wave equations we mean a solution of the nonlinear operator equation

$$(3.2) Au - Nu = f, u \in D(A), f \in H$$

where $Au = \sum_{j,k} u_{j,k} (\mu_j - \tau^2 k^2) \psi_j(x) e^{i\tau kt}$ for

$$u \in D(A) = \left\{ u = \sum_{j,k} u_{j,k} \psi_j(x) e^{i\tau kt} \, \left| \, \sum_{j,k} |u_{jk} (\mu_j - \tau^2 k^2)|^2 < \infty \right. \right\}$$

and Nu = F(t, x, u) for $u \in L_2(\Omega, \mathbb{R}^m)$.

Regarding $F = F_1 + F_2$, we assume

(3.3) F_1 is a Caratheodory function such that for some $a_1 > 0$, $k \in (0,1)$ and $h_1 \in L_2(\Omega, R)$ it satisfies

$$|F_1(t, x, y)| \le a_1 |y|^k + h_1(t, x)$$
 for a.e. $(t, x) \in \Omega$, all $y \in \mathbb{R}^m$.

(3.4) F_2 is a Caratheodory function and there are $h_2 \in L_2(\Omega, R)$ and some consecutive eigenvalues $\lambda_i < \lambda_{i+1}$ of A, and $0 < a_2 < \min\{\lambda - \lambda_i, \lambda_{i+1} - \lambda\}$ such that

$$|F_2(t, x, y) - \lambda y| \le a_2|y| + h_2(t, x)$$
, for a.e. $(t, x) \in \Omega$, $y \in \mathbb{R}^m$.

Our first nonresonance result for (3.1) is for monotone nonlinearities F and is an application of Theorem 2.1 (cf. [Mi-8]).

THEOREM 3.1 Let $Q=(0,\pi)$, $\mu_j\geq 0$ for each $j\in J$, and each nonzero eigenvalue of A is of finite multiplicity. Suppose that (3.3)–(3.4) hold and F is increasing in y. Then there is a T-periodic weak solution $u\in L_2$ of (3.1) for each $f\in L_2(\Omega,R^m)$.

Proof. Let $\{H_n\}$ be an increasing sequence of finite dimensional subspaces spanned by the eigenfunctions $\{\psi_j(x)e^{i\tau kt}\}$ and $P_n:H\to H_n$ be the orthogonal projections onto H_n . Since the eigenfunctions are dense in $L_2(\Omega,R^m)$, $\Gamma=\{H_n,P_n\}$ is a projection scheme for L_2 . Moreover, since $Nu=F(t,x,u):L_2\to L_2$ is monotone, the map $A-N:D(A)\subset L_2(\Omega,R^m)\to L_2(\Omega,R^m)$ is pseudo A-proper w.r.t. Γ . In view of Theorem 2.1 and Remark 2.2, it remains to verify conditions (2.2)-(2.3) with $X=H=L_2$.

First, note that the spectral gap of $A_{\lambda} = A - \lambda I$ induced by the gap $(\lambda_i, \lambda_{i+1})$ is $(\lambda_i - \lambda, \lambda_{i+1} - \lambda)$. Hence, A_{λ} has a bounded selfadjoint inverse in L_2 with the spectrum lying in $[(\lambda_i - \lambda)^{-1}, (\lambda_{i+1} - \lambda)^{-1}]$ and so $||A_{\lambda}^{-1}|| \leq \max\{(\lambda - \lambda_i)^{-1}, (\lambda_{i+1} - \lambda)^{-1}\} \leq 1/a$. Hence, $0 < a < ||A_{\lambda}^{-1}|| = \min\{|\mu| | \mu \in \sigma(A - \lambda I)\}$. Next, set $N_i u = F_i(t, x, u)$ for $u \in L_2$, i = 1, 2. Then, by the Minkowski and Hölder inequalities, (3.3)–(3.4) imply that for some constants c_1 and c_2

$$||N_1 u|| \le c_1 ||u||^k + ||h_1||$$
 and $||N_2 u|| \le c_2 ||u|| + ||h_2||$ for $u \in L_2$.

Since k < 1, there are positive a, b and r such that $a < \min\{\lambda - \lambda_i, \lambda_{i+1} - \lambda\}$ and

$$||Nu - \lambda u|| \le a||u|| + b$$
 for all $||u|| \ge r$.

Hence, (2.2)–(2.3) hold and Theorem 2.1 is applicable by Remark (2.2). \square

It is easy to see that condition (3.4) is implied by the following two conditions with $\lambda = (\lambda_i + \lambda_{i+1})/2$ and $a = (\lambda_{i+1} - \lambda_i)/2 - \epsilon$:

(3.5) There are constants M>0 and r>0 and $h\in L_2(\Omega,R)$ such that for each 1< l < m

$$|F_{2,l}(t,x,y)| \le M|y_l| + h_2(t,x)$$
 for a.e. $(t,x) \in \Omega, y \in \mathbb{R}^m$.

(3.6) For a.e.
$$(t, x) \in \Omega$$
, $y = (y_1, \dots, y_m) \in \mathbb{R}^m$ with $|y_l| \ge r$

$$\lambda_i + \epsilon \le F_{2,l}(t, x, y)/y_l \le \lambda_{i+1} - \epsilon.$$

Hence, when m = 1, we have

COROLLARY 3.1 Let $F: R \times (0, \pi) \times R \to R$ be a 2π -periodic in t Caratheodory function satisfying conditions (3.5)-(3.6) and $F(t, x, \cdot)$ be increasing. Then, for each $f \in L_2$, there is a 2π -periodic weak solution $u \in L_2$ of

(3.7)
$$u_{tt} - u_{xx} - F(t, x, u) = f(t, x), \qquad t \in R, \ x \in (0, \pi)$$
$$u(t, 0) = u(t, \pi) = 0, \qquad t \in R$$
$$u(t + 2\pi, x) = u(t, x), \qquad t \in R, \ x \in (0, \pi).$$

Remark 3.1 Without the monotonicity of F, the solvability of (3.7) for a dense set of f's in L_2 was proved by Hofer [Ho] under a global Lipschitz condition on F (cf. also [W]) and by Tanaka [Ta-1] without this condition. When F is monotone, Corollary 3.1 is due to Mawhin [Ma-2].

We continue our study of (3.1) when a nonlinear perturbation F satisfies asymptotic nonuniform nonresonance conditions with respect to two consecutive eigenvalues of the associated linear problem. These conditions are more general than (3.4) and (3.5)–(3.6) and our method of study requires a different approach based on Theorem 2.4. When m = n = 1, this problem has been studied by Mawhin-Ward [Ma-Wa] for the wave equation (3.7).

Let A be the abstract realization of the linear problem associated with (3.1) and $\lambda_i < \lambda_{i+1}$ be two consecutive eigenvalues of A having finite multiplicities.

Let $F: \Omega \times \mathbb{R}^m \to \mathbb{R}^m$ be a Caratheodory function such that for each r > 0 and $1 \le l \le m$ there are functions $a_l, b_l \in L_{\infty}(\Omega)$ and $h_r \in L_2(\Omega)$ such that

(3.8)
$$F(t,x,y)| \le h_r(t,x) for a.e. (t,x) \in \Omega, |y| \le r,$$

(3.9)
$$a_l(t,x) \le \liminf_{|y_l| \to \infty} F_l(t,x,y) y_l^{-1} \le \liminf_{|y_l| \to \infty} F_l(t,x,y) y_l^{-1} \le \beta_l(t,x)$$

uniformly a.e. in $(t, x) \in \Omega$ and $(y_1, \dots, y_{l-1}, y_{l+1}, \dots, y_m) \in \mathbb{R}^{m-1}$ and

(3.10)
$$\lambda_i \leq a_l(t,x) \leq \beta_l(t,x) \leq \lambda_{i+1} \quad \text{a.e. on } \Omega$$

with $\lambda_i < a_i(t, x)$ and $\beta_i(t, x) < \lambda_{i+1}$ on some sets of positive measure.

THEOREM 3.2 ([Mi-8]) Let (3.8)-(3.10) hold and

- (a) $\operatorname{sign} \lambda_i F$ be monotone if n = 1, i.e., $\operatorname{sign} \lambda_i (F(t, x, y) F(t, x, z)) \cdot (y z) \ge 0$ for a.e. $(t, x) \in \Omega$ and all $y, z \in \mathbb{R}^m$;
- (b) If n > 1, then $F = F_1 + F_2$ and for some positive constants c, k_1 , k_2 with $(k_1 + k_2)||A^{-1}|| < 1$, $k_2 < c$;

(3.11)
$$\operatorname{sign} \lambda_{i} (F_{1}(t, x, y) - F_{1}(t, x, z), y - z) \geq c|y - z|^{2}$$
for a.e. $(t, x) \in \Omega$, $y, z \in \mathbb{R}^{m}$,

(3.12)
$$|F_{i}(t, x, y) - F_{i}(t, x, z)| \le k_{i}|y - z|$$
for a.e. $(t, x) \in \Omega, y, z \in \mathbb{R}^{m} \quad i = 1, 2.$

Then there is a T-periodic weak solution $u \in L_2(\Omega, \mathbb{R}^m)$ of (3.1) for each $f \in L_2(\Omega, \mathbb{R}^m)$.

Proof. We have that either $0 < \lambda_i < \lambda_{i+1}$ or $\lambda_i < \lambda_{i+1} < 0$. We may assume that $\lambda_i > 0$, for otherwise instead of the corresponding operator equation

(3.13)
$$Au - Nu = f, u \in D(A), f \in H = L_2(\Omega, \mathbb{R}^m)$$

where Nu = F(t, x, u), we can consider the equivalent equation

$$A_1u - N_1u = -f$$

with $A_1 = -A$, $N_1 = -N$, $\sigma(A_1) = \{ \dots < 0 < -\lambda_{i+1} < -\lambda_i < \dots \}$. Then, setting

$$a_{1l}=-\beta_l,\beta_{1l}=-a_l\ \tilde{F}=-F,$$

we see that conditions (3.9), (3.10) and (3.12) hold with a_l , β_l , F, λ_i and λ_{i+1} replaced respectively by a_{1l} , β_{1l} , \tilde{F} , $-\lambda_i$ and $-\lambda_{i+1}$ and the function $\operatorname{sign}(-\lambda_i)\tilde{F}_1 = \operatorname{sign} \lambda_i F_1$ is monotone if n=1, or satisfies (3.11) if n>1. Hence, we can assume that $\lambda_i>0$ and therefore $N:H\to H$ is monotone when n=1 and N is c-strongly monotone and N_j are k_j -contractive, j=1,2, when n>1.

Next, we shall show that N is a bounded asymptotically $\{B_1, B_2\}$ -quasilinear map with $B_1 = C_1 - \epsilon I$ and $B_2 = C_2 + \epsilon I$ for some $\epsilon > 0$, where C_1 and C_2 are $m \times m$ diagonal matrices with the diagonal entries $a_1(t,x), \cdots, a_m(t,x)$ and $\beta_1(t,x), \cdots, \beta_m(t,x)$ respectively. By (3.9), for $\epsilon > 0$ there is an r > 0 such that for each $1 \leq l \leq m$, for $a.e.(t,x) \in \Omega$ and all $y = (y_1, \ldots, y_l, \ldots, y_m) \in R^m$ with $|y_l| \geq r$

(3.14)
$$a_l(t,x) - \epsilon \leq F_l(t,x,y)y_l^{-1} \leq \beta_l(t,x) + \epsilon.$$

Hence, by (3.8),

$$|F_l(t,x,y)| \le (\lambda_{i+1} + \epsilon)|y_l| + h_r(t,x)$$
 for a.e. $(t,x) \in \Omega$, $y \in \mathbb{R}^m$,

and therefore N is continuous and bounded in H. For each $1 \leq l \leq m$, define a function $G_l: \Omega \times \mathbb{R}^m \to \mathbb{R}$ by (cf. [Ma-Wa])

$$G_{l}(t,x,y) = \begin{cases} y_{l}^{-1}F_{l}(t,x,y), & \text{if } |y_{l}| \geq r \\ r^{-1}F_{l}(t,x,y_{1},\ldots,y_{l-1},r,y_{l+1},\ldots,y_{m})(y_{l}/r) & \\ + (1-y_{l}/r)a_{l}(t,x), & \text{if } 0 \leq y_{l} \leq r \\ r^{-1}F_{l}(t,x,y_{1},\ldots,y_{l-1},-r,y_{l+1},\ldots,y_{m})(y_{l}/r) & \\ + (1+y_{l}/r)a_{l}(t,x), & \text{if } -r \leq y_{l} \leq 0. \end{cases}$$

Using (3.14), it is easy to check that for $1 \le l \le m$

$$a_l(t,x) - \epsilon \le G_l(t,x,y) \le \beta_l(t,x) + \epsilon$$
 for a.e. $(t,x) \in \Omega, y \in \mathbb{R}^m$.

Moreover, $H_l(t, x, y) = F_l(t, x, y) - G_l(t, x, y)y_l$ is a Caratheodory function on $\Omega \times \mathbb{R}^m$ for $1 \leq l \leq m$ and

$$|H_l(t,x,y)| \leq 2h_r(t,x)$$
 for a.e. $(t,x) \in \Omega$, $y \in \mathbb{R}^m$.

Let G(t, x, y) be the $m \times m$ diagonal matrix with the diagonal entries

$$G_1(t, x, y), \ldots, G_m(t, x, y)$$
 for $(t, x) \in \Omega, y \in \mathbb{R}^m$.

Similarly, let N(t, x, y) be the $m \times m$ diagonal matrix with the diagonal entries $H_1(t, x, y), \ldots, H_m(t, x, y)$.

Now, for each $u \in H$, define $B(u): H \to H$ by B(u)v = G(t,x,u)v and Mu = H(t,x,u). Then, Nu = B(u)u + Mu on H, and it is easy to check that B(u) is selfadjoint and $B_1 \leq B(u) \leq B_2$ for each $u \in H$. Since $||Mu|| \leq 2||h_r||$ for each $u \in H$, it follows that N is asymptotically $\{B_1, B_2\}$ -quasilinear. Moreover, by Lemma 2.4, conditions (2.10)-(2.11) hold for any selfadjoint maps B^{\pm} with $C_1 \leq B^-$ and $B^+ \leq C_2$, where H^+ (resp., H^-) is a subspace of H spanned by the eigenfunctions of A corresponding to the eigenvalues $\lambda_k \leq \lambda_i$ (resp., $\lambda_k \geq \lambda_{i+1}$).

Finally, define the selfadjoint map $C_0: H \to H$ by $C_0u: a(t,x)u$, where a(t,x) is the $m \times m$ diagonal matrix with the diagonal entries $a_1(t,x), \ldots, a_m(t,x)$, for $(t,x) \in \Omega$. Since each $a_l(t,x) \geq \lambda_i > 0$ a.e. on Ω , it follows that C_0 is continuous and a_i -strongly monotone. Hence, $H(t,x) = A - (1-t)C_0 - tN$ is A-proper w.r.t. $\Gamma = \{H_k, P_k\}$ with $P_kAx = Ax$ on H_k for each $t \in [0,1)$ if $n \geq 1$, and H_1 is pseudo A-proper if n = 1 and is A-proper if n > 1 by Propositions 2.3 and 2.6 in [Mi-8]. Choosing $\epsilon > 0$ as in Lemma 2.4, we see that the conclusion of the theorem follows from Theorem 2.4. \square

Remark 3.2 When n=m=1, Theorem 3.2 with the Dirichlet boundary conditions for the semilinear wave equation (3.7) was proved by Mawhin-Ward [Ma-Wa] using rather different arguments based on the compactness of A^{-1} and the coincidence degree theory of Mawhin. Hence, their method is not applicable when n>1 since A^{-1} is not compact due to the presence of nonzero eigenvalues of infinite multiplicity.

Remark 3.3 In view of Remark 2.3, Theorem 3.2 still remains valid if $\lambda_i = 0$ (resp., $\lambda_{i+1} = 0$) provided each a_i (resp., β_i) is a constant.

IV. Semilinear Equations at Resonance

In this section, we study Eq. (1.1) when dimker $A = \infty$ and the nonlinearity interacts in some sense with the spectrum of A. Using a direct method, we first prove an extension of a theorem of Brezis- Nirenberg [Br-Ni-1] to pseudo A-proper maps A - N. Then, we use a perturbation method to study Eq. (1.1) assuming A - N is a closed map. Moreover, when N is a gradient map, we also give a necessary and sufficient condition for the solvability of Eq. (1.1).

4.1 Direct Method. Assume that $A:D(A)\subset H\to H$ satisfies:

- (4.1) There are positive constants a_{\pm} and a_0 such that
 - (i) $-a_{+}^{-1}||Ax||^{2} \le (Ax, x) \le a_{-}^{-1}||Ax||^{2}$ for $x \in D(A)$,
 - (ii) $||x|| \le a_0 ||Ax||$ for $x \in \tilde{H} = \ker A^{\perp}$.

Let Int(D) denote the interior of D and conv D be the convex hull of D. We have [Mi-7]

THEOREM 4.1 Let a linear closed map $A: D(A) \subset H \to H$ satisfy (4.1) and $N: H \to H$ be such that $\pm A + N: D(A) \subset H \to H$ is pseudo A-proper w.r.t. $\Gamma = \{H_n, P_n\}$ with $P_n Ax = Ax$ on H_n . Suppose that

(4.2) There are $\gamma < a_{\pm}$ and $\tau < a_0^{-2}(\gamma^{-1} - a_{\pm}^{-1})$ such that for every $y \in H$ and every $\delta > 0$ there exist $c_i(y)$, i = 1, 2, and $k(\delta)$ such that for each $x \in H$

$$(Nx - Ny, x) \ge \gamma^{-1} ||Nx||^2 - c_1(y)||Nx|| - \tau ||x_1|| - c_2(y)|| (\delta ||x_0|| + k(\delta)).$$

Then $\operatorname{Int}(R(A)+\operatorname{conv} R(N))\subset R(\pm A+N)$. Moreover, if N is onto, so is $\pm A+N$.

Proof. We shall consider only +A+N, since the case -A+N can be done in a similar way. Let $f \in \text{Int}(R(A) + \text{conv } R(N))$ be fixed and assume that

$$Ax_n + (1 - t_n)P_nPx_n + t_nP_nNx_n = t_nP_nf$$

for some $x_n \in H_n$, $t_n \in (0,1]$ and each n, where $P: H \to \ker A$ is the orthogonal projection. We shall show that $||x_n|| < r$ for some r and all n.

For each $h \in H$ with sufficiently small norm, there are $v \in D(A)$, $w_i \in H$ and $\lambda_i \geq 0$ with $\sum_{i=1}^k \lambda_i = 1$ such that

$$f + h = Av + \sum_{i=1}^{k} \lambda_i Nw_i.$$

Substituting this in the above equation, we get

$$(1-t_n)P_nPx_n+t_nP_n\left(Nx_n-\sum_{i=1}^k\lambda_iNw_i\right)+t_nP_nh=t_nP_nAv-Ax_n.$$

Taking the scalar product with x_n , we find

$$(1-t_n)||Px_n||^2+t_n\sum_{i=1}^k\lambda_i(Nx_n-Nw_i,x_n)+t_n(h,x_n)=t_n(Av,x_n)-(Ax_n,x_n).$$

Hence, by (4.2), for each $\delta > 0$,

$$(1-t_n)||Px_n||^2 + t_n \sum_{i=1}^k \lambda_i (\gamma^{-1}||Nx_n||^2 - c_1(w_i)||Nx_n|| - \tau ||x_{1n}||^2 - c_2(w_i) (\delta ||x_{0n}|| + k(\delta))) + t_n(h, x_n)$$

$$\leq t_n ||Av|| ||x_{1n}|| + a_+^{-1} ||Ax_n||^2,$$

or, after dropping the first term,

$$(h, x_n) + \gamma^{-1} ||Nx_n||^2 \le c(h) ||Nx_n|| + \tau ||x_{1n}||^2 + ||Av|| ||x_{1n}|| + (t_n a_+)^{-1} ||Ax_n||^2 + c(h)(\delta ||x_{0n}|| + k(\delta))$$

for some constant c(h). But, since

$$Ax_n + t_n(I - P)P_nNx_n = t_n(I - P)P_nf$$

we get that

$$||Ax_n|| < t_n(||Nx_n|| + ||f||)$$
 and $||x_{1n}|| < a_0t_n(||Nx_n|| + ||f||)$.

Hence,

$$\begin{split} (h,x_n) &\leq a_+^{-1} \left(||Nx_n|| + ||f|| \right)^2 - \gamma^{-1} ||Nx_n||^2 + \tau a_0^2 \left(||Nx_n|| + ||f|| \right)^2 \\ &+ c ||Nx_n|| + a_0 ||Av|| \left(||Nx_n|| + ||f|| \right) + c(h) \left(\delta ||x_{0n}|| + k(\delta) \right) \\ &= \left(a_+^{-1} - \gamma^{-1} + \tau a_0^2 \right) ||Nx_n||^2 + \left(c(h) + 2a_+^{-1} ||f|| + 2\tau a_0^2 ||f|| + a_0 ||Av|| \right) ||Nx_n|| + c(h) \left(\delta ||x_{0n}|| + k(\delta) \right) + c_1(h), \end{split}$$

where $c_1(h) = (a_+^{-1} + \tau a_0^2) ||f||^2 + a_0 ||Av|| ||f||$. Since $\tau a_0^2 < \gamma^{-1} - a_+^{-1}$, it follows that, for each $\delta > 0$ and some $c_2(h)$ large enough,

$$(h, x_n) \le c(h) (\delta ||x_{0n}|| + k(\delta)) + c_1(h) \le c_2(h) (\delta ||x_n|| + k(\delta)).$$

Suppose that $||x_n|| \to \infty$ as $n \to \infty$, and for r > 0, set $w(r) = \inf_{\delta < 0} \{\delta r + k(\delta)\}$ so that $w(r)/r \to 0$ as $r \to \infty$. Since $(h, x_n) \le w(||x_n||)c_2(h)$, the uniform boundedness principle implies that $||x_n||/w(||x_n||) \le C$, a contradiction. Hence, $||x_n|| < r$ for some r > 0 and all n. Consequently, $H_n(t, x) = Ax + (1 - t)P_nPx + tP_nNx \ne tP_nf$ on $[0, 1] \times \partial B(0, r) \cap H_n$ for each n, and the Brouwer degree

$$\deg(A + P_n N, B(0, r) \cap H_n, P_n f) = \deg(A + P_n P, B(0, r) \cap H_n, 0) \neq 0.$$

Thus, $Ax_n + P_n Nx_n = P_n f$ for some $x \in B(0, r) \cap H_n$ and each n, and consequently, Ax + Nx = f for some $x \in D(A)$ by the pseudo A-properness of A + N. \square

Now, condition (4.2) implies that $\limsup_{\|x\|\to\infty}\|Nx\|/\|x\| \le \gamma$, i.e. N is a quasibounded map with the quasinorm $|N| \le \gamma$. As pointed out in [Br-Ni-1], if N is potential and quasibounded, then a condition of type (4.2) holds. Moreover, if a is the smallest positive constant such that $(Ax, x) \ge -a^{-1}\|Ax\|^2$ on D(A) and A is selfadjoint, then $Ax = \lambda x$ implies that $\lambda \ge -a$ and -a is an eigenvalue of A. Hence, roughly speaking, the conditions $0 < \gamma < a$ and $\limsup_{\|x\|\to\infty}\|Nx\|/\|x\| \le \gamma$ mean that the nonlinearity N asymptotically stays away from the nonzero eigenvalues of A.

Note that a_{\pm} in (4.1) exist since, by the boundedness of $A^{-1}:R(A)\subset H\to H$, $||Ax_1||\geq a||x_1||$ for some a>0 and all $x_1\in D(A)\cap R(A)$, and therefore $(Ax,x)=(Ax,x_1)\geq -||Ax||\,||x_1||\geq -a^{-1}||Ax||^2$ for all $x\in D(A)$. Here we used the fact that $R(A)=N(A)^{\perp}$. If A is sefadjoint and $a<\infty$ is the smallest positive constant such that $(Ax,x)\geq -a^{-1}||Ax||^2$ on D(A), we have shown above that -a is the largest eigenvalue of A less than 0. More generally, suppose that $A:D(A)\subset H\to H$ is a normal linear map with closed range. If H_c is the complexification of the real Hilbert space H and $A_c:D(A_c)\subset H_c\to H_c$ is defined by $A_c(x+iy)=Ax+iAy$ on $D(A_c)=\{x+iy\mid x,y\in D(A)\}$, then Hetzer [He] has shown that

$$a = a(A) = \inf\{|\lambda|^2/(-\operatorname{Re}\lambda) \mid \lambda \in \sigma(A_c), \operatorname{Re}\lambda < 0\}.$$

In general a could be infinite, in which case A is a linear monotone map, and could belong to $\sigma(A_c)$ or be a regular value of A_c .

Theorem 4.1 has been proved by Brezis-Nirenberg [Br-Ni-1] when $N: H \to H$ is a monotone map and the partial inverse A^{-1} is compact using the Leray-Shauder and monotone operator theories. In view of the various examples of pseudo A-proper maps A+N discussed in [Mi-4-8], this result holds also for many other classes of maps A and N, even when A^{-1} is not compact. For example, we have the following corollary.

COROLLARY 4.1 Let A and N satisfy conditions (4.1)-(4.2) and $N = N_1 + N_2$ be such that N_1 is c-strongly monotone, k_1 -ball contractive and N_2 is k_2 ball contractive, with k_1 , k_2 sufficiently small, and continuous. Then A+N is surjective.

Proof. The map $A + N : D(A) \subset H \to H$ is A-proper by Proposition 2.6 in [Mi-8] (cf. also [Mi-4]). Moreover, since N is c-strongly monotone, it is well-known to be subjective and the conclusion follows from Theorem 4.1. \square

For our second corollary, we need to recall the following result proven in [Mi-7].

PROPOSITION 4.1 Let $A:D(A)\subset H\to H$ be selfadjoint, H^{\pm} be closed subspaces of H with $H=H^-\oplus H^+$ and $H^-\cap H^+=\{0\}$ and $\Gamma=\{H_n,P_n\}$ be a scheme for H that satisfies (2.21). Suppose that $N:H\to H$ has a symmetric weak Gateaux derivative N'(x) on H and there are symmetric maps $B^{\pm}\in L(H)$ such that $B^-\leq N'(x)\leq B^+$ for each $x\in H$ and

(4.3)
$$((A - B^{-})x, x) \le 0 \quad \text{for } x \in D(A) \cap H^{-};$$
(4.4)
$$((A - B^{+})x, x) > 0 \quad \text{for } x \in D(A) \cap H^{+}.$$

$$(4.4) ((A - B^{+})x, x) \ge 0 for x \in D(A) \cap H^{+}.$$

Then $\pm A + N : D(A) \subset H \to H$ is pseudo A-proper w.r.t. Γ .

In view of Proposition 4.1, we have the following special case of Theorem 4.1.

COROLLARY 4.2 Let $A: D(A) \subset H \to H$ and $N: H \to H$ be as in Proposition 4.1 and satisfy conditions (4.1)-(4.2). Then $\operatorname{Int}(R(A) + \operatorname{conv} R(A)) \subset R(\pm A + N)$ and $\pm A + N$ is surjective if such is N.

Next, we shall discuss some conditions that imply (4.3)–(4.4). Let $A:D(A)\subset$ $H \to H$ be selfadjoint and suppose that

- (4.5) $B^{\pm} = \sum_{i=1}^{m} \lambda_{i}^{\pm} P_{i}^{\pm}$ commute with A, where $P_{i}^{\pm} : H \to \ker(B^{\pm} \lambda_{i}^{\pm} I)$ are orthogonal projections, $\lambda_{1}^{\pm} \le \cdots \le \lambda_{m}^{\pm}$ and are pairwise disjoint;
- (4.6) $\bigcup_{i=1}^{m} [\lambda_i^-, \lambda_i^+] \subset \overline{\rho(A)}$, the closure of the resolvent set of A.

As we have seen in Section 2.2, $A - B^{\pm}$ have the spectral resolution

$$A - B^{\pm} = \sum_{i=1}^{m} \int_{-\infty}^{\infty} (\lambda - \lambda_i^{\pm}) dE_{\lambda} P_i^{\pm}.$$

Let δ be small enough such that $\mu_i^{\pm} = \lambda_1^{\pm} \mp \delta$ satisfy

(4.7)
$$\bigcup_{i=1}^{m} [\mu_i^-, \mu_i^+] \subset \rho(A).$$

Then the operators $B_{\delta}^{\pm} = B^{\pm} \mp \delta I$ have μ_{i}^{\pm} as their eigenvalues and $\ker(B_{\delta}^{\pm} - \mu_{i}^{\pm}) = \ker(B^{\pm} - \lambda_{i}^{\pm})$. Since B_{δ}^{\pm} commute with A, the spectral resolutions of $A - B_{\delta}^{\pm}$ are

(4.8)
$$A - B^{\pm} = \sum_{i=1}^{m} \int_{-\infty}^{\infty} (\lambda - \mu_i^{\pm}) dE_{\lambda} P_i^{\pm}.$$

Define the orthogonal projections P^{\pm} by

$$P^{+} = \sum_{i=1}^{m} E(\mu_{i}^{+}, \infty) P_{i}^{+} \quad \text{and} \quad P^{-} = \sum_{i=1}^{m} E(-\infty, \mu_{i}^{-}) P_{i}^{-}$$

and let $H^{\pm} = P^{\pm}(H)$. Note that by (4.7)

(4.9)
$$P^{+} = \sum_{i=1}^{m} E(\mu_{i}^{+}, \infty) P_{i}^{+}$$

and

$$\operatorname{dist}\bigg(\bigcup_{i=1}^m [\mu_i^-, \mu_i^+], \sigma(A)\bigg) \geq \delta.$$

Moreover, it follows from (4.8) that $((A - B_{\delta}^-)x, x) \leq -\delta ||x||^2$ on $D(A) \cap H^-$ and $((A - B_{\delta}^-)x, x) \geq \delta ||x||^2$ on $D(A) \cap H^+$. Hence, conditions (4.3)–(4.4) hold. Moreover, if we assume that $P_i^- = P_i^+$ for $1 \leq i \leq m$, then by (4.9), $P^+ = I - P^-$, i.e., $H^+ = (H^-)^{\perp}$. \square

By the discussion above, we have

COROLLARY 4.3 Let $A: D(A) \subset H \to H$ be selfadjoint, H^{\pm} be closed subspaces of H with $H = H^{-} \oplus H^{+}$ and $H^{-} \cap H^{+} = \{0\}$ and $\Gamma = \{H_{n}, P_{n}\}$ satisfy (2.21). Suppose that $N: H \to H$ has a symmetric weak Gateaux derivative N'(x) on H and $B^{-} \leq N'(x) \leq B^{+}$ for $x \in H$ and some selfadjoint maps $B^{\pm} \in L(H)$. If conditions (4.1)-(4.2) and (4.5)-(4.6) hold, then Int $(R(A) + \operatorname{conv} R(A)) \subset R(\pm A + N)$ and $\pm A + N$ is surjective if such is N.

A particular case of Corollary 4.3 is when $B^- = \lambda_k I$ and $B^+ = \lambda_{k+1} I$ for two consective eigenvalues $\lambda_k < \lambda_{k+1}$ of A. If H^- (resp. H^+) is the subspace of H spanned by the eigenvectors of A corresponding to the eigenvalues $\lambda_i \leq \lambda_k$ (resp. $\lambda_i \geq \lambda_{k+1}$), then $H^+ = (H^-)^{\perp}$ and $H = H^- \oplus H^+$. Let $\Gamma = \{H_n, P_n\}$ be a scheme for H with $P_n Ax = Ax$ and $P^{\pm} : H \to H^{\pm}$ be the orthogonal projections onto H^{\pm} . Then Γ satisfies (2.21) with $H_n = H_n^- \oplus H_n^-$ and $H_n^{\pm} = H_n \cap H^{\pm}$.

Now, since $B^- \leq N'(x) \leq B^+$ on H, the mean value theorem implies that $||Nx - Ny|| \leq \max\{||B^-||, ||B^+||\} ||x - y||$ for $x, y \in H$. In particular, taking $B^- = \lambda_k I$ and $B^+ = \lambda_{k+1} I$, we get

$$\limsup_{\|x\|\to\infty} \|Nx\|/\|x\| \le \max\left\{|\lambda_k|, |\lambda_{k+1}|\right\},\,$$

and, by (4.2) as noted above

$$\limsup_{\|x\|\to\infty} \|Nx\|/\|x\| \le \gamma < a,$$

Hence, for N to interact with λ_k (or λ_{k+1}), it must be the zero eigenvalue.

4.2 A Perturbation Method. In this section we shall present another way of studying Eq. (1.1) by looking at the perturbed equations

$$(4.10) Ax + Nx + \epsilon Gx = f, \quad \epsilon > 0$$

where G is a bounded map. This approach consists of three basic steps. The first step is to establish the solvability of (4.10) for each $\epsilon \in (0, \epsilon_0)$; the second step is to obtain an a priori bound on the solutions, i.e. to show that the set $\{x_{\epsilon} \mid x_{\epsilon} \text{ is a solution of } (4.10)\}$ is bounded in a suitable space as $\epsilon \to 0$, and the final step is to show that a weak limit of $\{x_{\epsilon}\}$ is a solution of Ax + Nx = f.

The approach requires a closedness property of A + N, which is defined next.

Definition 4.1 Let X be a Banach space embedded in H. We say that $A+N: D(A)\cap X \to H$ satisfies condition (*) if whenever $\{x_{\epsilon} \mid Ax_{\epsilon}+Nx_{\epsilon}+\epsilon Gx_{\epsilon}=$

= f, $0 < \epsilon < \epsilon_0$ is bounded as $\epsilon \to 0$ (in X or H), then there is an $x \in D(A)$ such that Ax + Nx = f.

The following result gives various conditions on A and N which guarantee that condition (*) holds.

PROPOSITION 4.2 Let $A: D(A) \cap X \to H$ be a closed linear densely defined map and N be a nonlinear map. Suppose that either A^{-1} is compact and N monotone, or the conditions of Corollary 4.1 and Proposition 4.1 hold. Then A+N satisfies condition (*).

Proof. (Sketch) Let $\{x_{\epsilon} \mid Ax_{\epsilon}+Nx_{\epsilon}+\epsilon Gx_{\epsilon}=f, \ 0<\epsilon<\epsilon_{0}\}$ be bounded in a corresponding space. Then, using similar arguments as in the proofs of Propositions 5.1 in [Mi-4], Corollary 4.1 and Proposition 4.1, one can show that the weak limit x of $\{x_{\epsilon}\}$ solves the equation Ax + Nx = f. \square

The problem of getting a priori estimates for (4.10), i.e. of showing the boundedness of $\{x_{\ell}\}$, is harder to handle. Our basic result is [Mi-6]:

THEOREM 4.2 Let $A: D(A) \subset H \to H$ be a linear densely defined selfadjoint map with $R(A) = N(A)^{\perp}$ and a > 0 be the largest number such that $(Ax, x) \geq -a^{-1}||Ax||^2$ on D(A). Suppose that $N: H \to H$ is nonlinear, condition (*) holds and

(4.11) There is a decomposition $f = f^* + f^{**}$ with $f^* \in R(A)$ such that for some $\gamma < a$ and a constant c

$$(Nx - f^{**}, x) \ge \gamma^{-1} ||Nx||^2 - c \text{ for } x \in H,$$

- (4.12) Eq. (4.10) with G = I is solvable for each $0 < \epsilon < \epsilon_0$ and either
- (4.13) the set $\{x_{\epsilon} \mid Ax_{\epsilon} + Nx_{\epsilon} + \epsilon x_{\epsilon} = f, ||Ax_{\epsilon}|| \le C, ||Nx_{\epsilon}|| \le C \text{ for all } 0 < \epsilon < \epsilon_0 \text{ and some } C\}$ is bounded in H, or
- (4.14) $||N(x_0+x_1)|| \to \infty$ as $||x_0|| \to \infty$, $x_0 \in N(A)$, uniformly for x_1 in bounded subsets of R(A).

Then Eq. (1.1) is solvable.

Proof. Let $\epsilon > 0$ be small and $Ax_{\epsilon} + Nx_{\epsilon} + \epsilon x_{\epsilon} = f$ for some $x_{\epsilon} \in D(A)$. We need to show that $\{x_{\epsilon}\}$ is bounded as $\epsilon \to 0$. Taking the inner product with x_{ϵ} , we find

$$\epsilon ||x_{\epsilon}||^2 + (Nx_{\epsilon} - f, x_{\epsilon}) = -(Ax_{\epsilon}, x_{\epsilon}) \le a^{-1} ||Ax_{\epsilon}||^2,$$

and by (4.11),

$$|\epsilon||x_{\epsilon}||^{2} + \gamma^{-1}||Nx_{\epsilon}||^{2} \le -(Ax_{\epsilon}, x_{\epsilon}) + (f^{**}, x_{\epsilon}) + C||Nx_{\epsilon}|| + C.$$

Since $A\nu = f^*$ for some $\nu \in D(A)$ and $(A(x_{\epsilon} - \nu), x_{\epsilon} - \nu) \ge -a^{-1} ||A(x_{\epsilon} - \nu)||^2$, we get

$$\frac{\epsilon ||x_{\epsilon}||^{2} + \gamma^{-1}||Nx_{\epsilon}||^{2}}{\leq a^{-1}||A(x_{\epsilon} - \nu)||^{2} - (A(x_{\epsilon} - \nu), \nu) + C||Nx_{\epsilon}|| + C}{\leq a^{-1}||Ax_{\epsilon}||^{2} + C||Ax_{\epsilon}|| + C||Nx_{\epsilon}|| + C}.$$

For each $\gamma < \beta < a$ there is a $C(\beta)$ such that

$$\gamma^{-1}||Nx_{\epsilon}||^{2} - C||Nx_{\epsilon}|| \ge \beta^{-1}||Nx_{\epsilon}||^{2} - C(\beta).$$

Similarly, for each $\beta < \delta < a$ there is a $C(\delta)$ such that

$$\epsilon ||x_{\epsilon}||^{2} + \beta^{-1}||Nx_{\epsilon}|| \leq a^{-1}||Ax_{\epsilon}||^{2} + C||Ax_{\epsilon}|| + C(\beta) \leq \delta^{-1}||Ax_{\epsilon}||^{2} + C(\beta, \delta),$$

where $C(\beta, \delta) = C(\beta) + C(\delta)$. Since $Nx_{\epsilon} = f - Ax_{\epsilon} - \epsilon x_{\epsilon}$, the last inequality implies that

$$(4.15) \quad \epsilon \left[(1 + \epsilon \beta^{-1}) \|x_{\epsilon}\|^{2} - 2\beta^{-1} \|f\| \|x_{\epsilon}\| \right] \\ + \left(\beta^{-1} - \delta^{-1} - 2\epsilon(a\delta)^{-1} \right) \|Ax_{\epsilon}\|^{2} - 2\delta^{-1} \|f\| \|Ax_{\epsilon}\| \le C(\beta, \delta).$$

This implies that for some sufficiently small constants $C_i > 0$, i = 1, 2, independent of ϵ ,

$$\epsilon C_1 ||x_{\epsilon}||^2 + C_2 ||Ax_{\epsilon}||^2 \le C_3,$$

and therefore $||Ax_{\epsilon}|| \leq C$ for all small ϵ . Moreover, $||Nx_{\epsilon}|| \leq C$ by (4.15) and consequently $||x_{\epsilon}|| \leq C$ for all small ϵ if A + N satisfies (4.13). Thus, $Ax_{\epsilon} + Nx_{\epsilon} = f - \epsilon x_{\epsilon} \to f$ as $\epsilon \to 0$ and Ax + Nx = f for some $x \in D(A)$ by condition (*).

Next, if instead of (4.13) we assume (4.14), then $x_{\epsilon} = x_{0\epsilon} + x_{1\epsilon}$ and $||x_{1\epsilon}|| = ||A^{-1}Ax_{\epsilon}|| \leq C||A^{-1}||$ for all small ϵ . Moreover, $||x_{0\epsilon}|| \leq C$ for such ϵ by (4.14), and consequently $\{x_{\epsilon}\}$ is bounded. Then, the conclusion follows as above. \square

When N is a gradient map, we have the following sufficient and necessary conditions for the solvability of (1.1) (see [Mi-6]).

THEOREM 4.3 Let $A: D(A) \subset H \to H$ be a linear selfadjoint map and a > 0 be the largest number such that $(Ax, x) \ge -a^{-1}||Ax||^2$ on D(A). Assume that either (4.13) or (4.14) holds, $N = \partial F$ for some convex function $F: H \to R$ and

(4.16)
$$\limsup ||Nx||/||x|| < a/2, \quad as \ ||x|| \to \infty.$$

Suppose that either one of the following conditions holds

- (4.17) $A^{-1}: R(A) \subset H \to H$ is compact,
- (4.18) $0 \in \sigma(A)$ and $\sigma(A) \cap (0, \infty) \neq \emptyset$ and consists of eigenvalues of finite multiplicities,
- (4.19) There are closed subspaces H^{\pm} of H with $H = H^{-} \oplus H^{+}$ and $H^{-} \cap H^{+} = \{0\}$, a scheme $\Gamma = \{H_{n}, P_{n}\}$ satisfying (2.21) and N has a symmetric weak Gateaux derivative N' on H such that for some symmetric maps $B^{\pm} \in L(H)$ with $B^{-} \leq N'(x) \leq B^{+}$ on H we have
 - (i) $((A B^-)x, x) \le 0$ for $x \in D(A) \cap H^-$;
 - (ii) $((A-B^+)x, x) \ge 0$ for $x \in D(A) \cap H^+$

Then, Eq. (1.1) is solvable if and only if $f = f^* + f^{**}$ with $f^* \in R(A)$ and $f^{**} \in R(N)$.

Proof. If x is a solution of Ax + Nx = f, then $f^* = Ax$ and $f^{**} = Nx$. Conversely, let $f = f^* + f^{**}$ with $f^* \in R(A)$ and $f^{**} \in R(N)$. Then (4.16) and Proposition A.4 in [Br-Ni-1] imply that N and $N_{\epsilon} = N + \epsilon I$ satisfy (4.2) and (4.11), respectively, for each $\epsilon > 0$ small. Moreover, N_{ϵ} is monotone and $||N_{\epsilon}x|| \to \infty$ as $||x|| \to \infty$ and therefore it is surjective.

Since $A + N_{\epsilon}$ is pseudo A-proper by Propositions 2.2 in [Mi-4] and 4.1, Eq. (1.1) is solvable for each small $\epsilon > 0$ by Theorem 4.1. Finally, since condition (*) holds by Proposition 4.2, Eq. (1.1) is solvable by Theorem 4.2. \square

When $H = L_2(Q, \mathbb{R}^m)$ we can relax condition (4.14) in Theorem 4.2 (cf. [Mi-6]).

COROLLARY 4.4 Let $Q \subset R^n$ be a bounded domain, $H = L_2(Q, R^m)$ and $Nu = D_u F(x, u)$ with $F: Q \times R^m \to R$ measurable in x and convex and C^1 in u. Then, condition (4.14) in Theorem 4.2 can be weakened to

$$\int_{Q} F(x, u_0(x)) dx \to \infty \quad as \ ||u_0||_{L_2} \to \infty, \ u_0 \in N(A).$$

Remark 4.2 When n = 1, a/2 in (4.16) can be replaced by a by Proposition A.6 in [Br-Ni-1].

V. Periodic Solutions of Semilinear Wave Equations at Resonance

In this section we shall apply some of the abstract results from the previous section to semilinear (systems) of wave equations (3.1) allowing some type of interaction of the nonlinearity F with the spectrum of the linear problem.

We begin with strongly monotone and contractive nonlinearities. Suppose $\Omega = (0,T) \times Q$ and

(5.1) Let $F_i: \Omega \times \mathbb{R}^m \to \mathbb{R}^m$ be Caratheory functions, i = 1, 2, such that for some positive constants k_i and c, with k_1, k_2 sufficiently small,

$$(F_1(t, x, y) - F_1(t, x, z)) \cdot (y - z) > c|y - z|^2$$

and

$$|F_i(t,x,y) - F_i(t,x,z)| \le k_i |y-z|$$

for a.e. $(t, x) \in \Omega$ and all $y, z \in \mathbb{R}^m$.

(5.2) There are a>0 and $h\in L_2(\Omega)$ such that for $\gamma_1=1,$ and $\gamma_2\in(0,1)$

$$|F_i(t,x,y)| \le a|y|^{\gamma_i} + h(t,x) \quad \text{for a.e. } (t,x) \in \Omega, \ y \in R^m.$$

Let L and V be as introduced in Section III. Then, we have shown there that a linear map $A: D(A) \subset L_2 \to L_2$, induced by (3.1), is selfadjoint and $R(A) = N(A)^{\perp}$. Moreover, condition (4.1) holds for some a_{\pm} and a_0 .

We have

THEOREM 5.1 Let (5.1)-(5.2) hold with $a < a_+/2$ and m > 1 and $F = \partial \psi$ for some function $\psi(t,x,y): \Omega \times R^m \to R$ measurable in (t,x) and differentiable and convex in y. Then there is a T-periodic weak solution $u \in L_2$ for each $f \in L_2$.

Proof. Define $N_i u = F_i(t, x, u)$ for $u \in L_2$ and let $N = N_1 + N_2$. By (5.1), N_1 is c-strongly monotone and k_1 -ball contractive and N_2 is k_2 -ball contractive. Hence, A + N is A-proper w.r.t. $\Gamma = \{H_n, P_n\}$ for L_2 with $P_n A u = A u$ for $u \in H_n$ by Proposition 2.6 in [Mi-8] (cf. also [Mi-4]). Moreover, $||N_2 u||/||u|| \to 0$ as $||u|| \to \infty$ and

$$(Nu, u) = (N_1u, u) - ||N_2u|| ||u|| \ge c||u||^2 - (||N_10|| + ||N_2u||) ||u||.$$

Thus, N is coercive, i.e. $(Nu, u)/||u|| \to \infty$ as $||u|| \to \infty$, and monotone and therefore, it is surjective.

Next, by Proposition A.4 in [Br-Ni-1], there is a $\gamma < a_+$ such that

$$(Nu - Nv, u) \ge \gamma^{-1} ||Nu||^2 - c(v)$$
 for all $u, v \in L_2$

and some constant c(v). Hence, by Theorem 4.1, the equation Au - Nu = f is solvable for each $f \in L_2$, i.e., there is a T-periodic weak solution $u \in L_2$ of (3.1) for each $f \in L_2$. \square

As it will be seen below, conditions on F can be greatly relaxed when n = 1. Next, we shall study the solvability of (3.1) for a given f.

THEOREM 5.2 Let $G(t,x,u): \Omega \times \mathbb{R}^m \to \mathbb{R}$ be T-periodic in t, measurable in (t,x) and convex and C^1 in u and $F(t,x,u) = D_uG(t,x,u)$ satisfy (5.2). Suppose that

(5.3)
$$\int_{\Omega} G(t,x,u)dtdx \to \infty \quad as \ ||u_0||_{L_2} \to \infty, \ u_0 \in N(A).$$

Moreover, if n > 1 assume that $a < a_{+}/2$, (4.6) holds and

- (5.4) $F \in C^1(\mathbb{R}^m, \mathbb{R}^m)$ and has a symmetric derivative $D_y F(t, x, y)$ for a.e. $(t, x) \in \Omega$ and all $y \in \mathbb{R}^m$,
- (5.5) there exist two commuting $m \times m$ symmetric matrices b^{\pm} such that $b^{-} \leq D_{y}F(t,x,y) \leq b^{+}$ for a.e. $(t,x) \in \Omega$, $y \in R^{m}$

and, if n = 1, assume that $a < a_+$.

Then (3.1) has a T-periodic weak solution $u \in L_2$ if and only if $f = f^* + f^{**}$ with $f^* \in R(A)$ and $f^{**} \in R(N)$, where Nu = F(t, x, u).

Proof. Let n=1 and $a < a_+$. Then $A^{-1}: R(A) \subset L_2 \to L_2$ is compact and N satisfies condition (4.2) by Proposition A.1 in [Br-Ni-1]. If n>1 and $a < a_+/2$, condition (4.2) holds by Proposition A.4 in [Br-Ni-1]. Let $\lambda_1^{\pm} \leq \cdots \leq \lambda_m^{\pm}$ be the eigenvalues of b^{\pm} and set $B^{\pm}u = b^{\pm}u$. Then $B^{\pm} = \sum_{i=1}^{m} \lambda_1^{\pm} P_i^{\pm}$, where P_i^{\pm} is the orthogonal projection onto $\ker(B^{\pm} - \lambda_i^{\pm}I)$. It is easy to show that, after a possible renumeration of the eigenvalues of b^{\pm} , we may assume that $P_i^{+} = P_i^{-}$ for each i. Then, by our discussion in Section 4.1 the conclusion of the theorem follows from Theorem 4.3 and Corollary 4.4. \square

COROLLARY 5.1 Let n = 1 and $F(t, x, u) = D_u G(t, x, u)$ satisfy (5.2) with $h \in L_{\infty}(\Omega)$, $a < a_+$, and

(5.6)
$$G(t, x, u) \to \infty$$
 as $|u| \to \infty$ uniformly a.e. in Ω .

Then there is 2π periodic weak solution $u \in L_2$ of

(5.7)
$$\begin{cases} u_{tt} + u_{xxx} + F(t, x, u) = f(t, x), & t \in R, x \in (0, \pi), \\ u(t, 0) = u(t, \pi) = 0, & t \in R \\ u_{xx}(t, 0) = u_{xx}(t, \pi) = 0 \\ u(0, x) - u(2\pi, x) = u_t(0, x) = u_t(2\pi, x) = 0, & t \in R, x \in (0, \pi) \end{cases}$$

if and only if $f = f^* + f^{**}$ with $f^* \in R(A)$ and $f^{**} \in R(N)$.

Proof. Since $h \in L_{\infty}$ and G is convex in u, as in [Ma-Wi] there is $\delta > 0$ such that

(5.8)
$$G(t,x,u) \ge \delta|u| - h(t,x) \text{ for a.e. } (t,x) \in \Omega, \ u \in \mathbb{R}^m.$$

Since the L_1 and L_2 norms are equivalent on N(A) (cf. [Ba-Sa]), this implies (5.3) and the conclusion follows from Theorem 5.2. \square

When f = 0, we refer to [Ma-Wi] for related results obtained by a dual variational method. When n = 1, our results imply the following ones of Bahri-Brezis [Ba-Br] and Bahri-Sanchez [Ba-Sa] for the wave equation and (5.7), respectively.

Corollary 5.2 Let $g:R\to R$ be continuous nondecreasing and for some a<3, b and $\delta>0$

(5.9)
$$|g(u)| \le a|u| + b$$
 for all $u \in R$,

(5.10)
$$f = f^* + f^{**} \in L_{\infty}(\Omega)$$
 with $f^* \in R(A)$ and $g(-\infty) + \delta \leq f^{**}(t,x) \leq g(+\infty) - \delta$ a.e. in Ω .

Then there is a weak solution of (5.7) with F replaced by g, and of

(5.11)
$$\begin{cases} u_{tt} - u_{xx} + g(u) = f(t, x), & t \in R, \ x \in (0, \pi), \\ u(t, 0) = u(t, \pi) = 0, & t \in R \\ u(t + 2\pi, x) = u(t, x), & t \in R, \ x \in (0, \pi). \end{cases}$$

Proof. Consider first (5.7). We note that it is equivalent to

(5.12)
$$Av + g(u^* + v) - f^{**} = Av + D_v G(t, x, v) = 0$$

where $u^* \in L_{\infty}(\Omega)$ is the solution of $Au = f^*$, $u = u^* + v$, and

$$G(t,x,v) = \int_0^v \left[g(u^* + s) - f^{**}(t,x) \right] ds.$$

Since (5.9)–(5.10) imply

$$(g(u) - f^{**}(t, x)) u \ge \delta |u|/2 - C \text{ for } (t, x) \in \Omega, \ u \in R,$$

we have that (cf. [Mi-6]) (5.8) and therefore, (5.3) holds. Hence, the weak solvability of (5.7) follows from Theorem 5.2 applied to (5.12).

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Next, consider (5.11). Conditions (5.9)–(5.10) and the monotonicity of g imply (4.11), (4.12) and (*) and by Theorem 4.1 and Proposition 4.2. As shown by the proof in $[\mathbf{Br}]$, condition (4.13) also holds, and consequently the weak solvability of (5.11) follows from Theorem 4.2. \square

Another condition on f^{**} is given in the following result (cf. [Mi-8]).

Theorem 5.3 Let the conditions of Theorem 5.2 hold. Suppose that $f = f^* + f^{**} \in L_2(\Omega)$ with $f^* \in R(A)$ and

$$(5.13) (F(t,x,y) - f^{**}(t,x)) \cdot y \ge c|F(t,x,y) - f^{**}(t,x)||y|$$

for all $|y| \ge R$, a.e. $(t, x) \in \Omega$, some R, and a sufficiently small c > 0. Then there is a T-periodic weak solution of (3.1).

Proof. We shall first show that (5.2) and (5.13) imply condition (4.11). We have

$$(Nu - f^{**}, u) = \int_{|u(t,x)| < R} (F(t, x, u(t,x)) - f^{**}(t,x)) \cdot u(t,x) dt dx$$

$$+ \int_{|u(t,x)| \ge R} (F(t, x, u(t,x)) - f^{**}(t,x)) \cdot u(t,x) dt dx$$

$$= I_1 + I_2.$$

By (5.2),

$$|I_1| \le \int_{|u(t,x)| < R} (a|u(t,x)| + h(t,x) + |f^{**}(t,x)|) |u(t,x)| dt dx$$

$$\le aR^2 + R(||h|| + ||f^{**}||) = c_1.$$

By (5.2) and (5.13), for $|u(t,x)| \ge R$

$$\begin{split} &(F(t,x,u(t,x)) - f^{**}(t,x)) \cdot u(t,x) \\ &\geq \frac{c}{a} \left| F(t,x,u(t,x)) - f^{**}(t,x) \right| \cdot (|F(t,x,u(t,x))| - h(t,x)) \\ &= \frac{c}{a} |F(t,x,u(t,x))|^2 - \frac{c}{a} |F(t,x,u(t,x))| \cdot (|f^{**}(t,x)| - h(t,x)) \\ &- \frac{c}{a} |f^{**}(t,x)| h(t,x). \end{split}$$

Hence,

$$I_2 \ge \frac{c}{a} ||Nu||^2 - \left(\frac{c}{a} ||f^{**}|| + ||h||\right) ||Nu|| - \frac{c}{a} ||f^{**}|| ||h||.$$

Since $C \le \epsilon C^2 + c(\epsilon)$ for each $\epsilon > 0$, condition (4.11) follows easily from the estimates on I_1 and I_2 .

Next, in view of Propositions A.1 and A.4 in [Br-Ni-1], condition (4.14) holds and therefore the equation $Ax + Nx + \epsilon x = f$ is solvable for each $\epsilon > 0$ as

shown in the proof of Theorem 4.3. Hence by our discussion in Section 4.2, the conclusion follows from Theorem 4.2.

Remark 5.1 It is easy to see that condition (5.10) implies condition (5.13).

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