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UNITED NATIONS EDUCATIONAL, SCIENTIFIC AND CULTURAL ORGANIZATION



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SMR 281/11

COLLEGE ON VARIATIONAL PROBLEMS IN ANALYSIS

(11 January - 5 February 1988)

AN INTRODUCTION TO CRITICAL POINT THEORY (PART II)

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II - Minimax theorems

Consider the following intuitive situation. If $g \in C^1(\mathbb{R}^2, \mathbb{R})$, we can view $g(x, y)$ as the altitude of the point of the graph of g having (x, y) as projection on $\mathbb{R}^2$.

Assume that there exists a bounded open neighborhood Ω of $u_0 \in \mathbb{R}^2$ and $u_1 \in \mathbb{R}^2 \setminus \Omega$ such that $g(u) > \max\{g(u_0), g(u_1)\}$ whenever $u \in \partial\Omega$. We can then consider the point $[u_0, g(u_0)]$ as located in a valley surrounded by a ring of mountains pictured by the set $\{[u, g(u)] : u \in \partial\Omega\}$, the point $[u_1, g(u_1)]$ being located outside the ring. To go from $[u_0, g(u_0)]$ to $[u_1, g(u_1)]$ in a way which minimizes the highest altitude of the path, we must cross the mountain ring through the lowest mountain pass. The projection on $\mathbb{R}^2$ of the top of this mountain pass will provide a critical point of g .

More generally, we have the following result.

Theorem 3. (Ambrosetti-Rabinowitz mountain pass theorem). Assume that $g \in C^1(X, \mathbb{R})$ satisfies (PS). If there exist $u_0 \in X$, $u_1 \in X \setminus \Omega$, where Ω is a bounded open neighborhood of u_0 , such that, for some $\alpha \in \mathbb{R}$,

$$g(u) \geq \alpha > \max\{g(u_0), g(u_1)\}$$
 whenever $u \in \partial\Omega$, then g has a critical value $c \geq \alpha$.

Proof. Let us define

$$\Gamma = \{ \gamma \in C([0,1], X) : \gamma(0) = u_0, \gamma(1) = u_1 \}$$

$$c = \inf_{\gamma \in \Gamma} \max_{t \in [0,1]} \varphi(\gamma(t)).$$

Since every curve $\gamma \in \Gamma$ crosses $\partial\Omega$, $\max_{t \in [0,1]} \varphi(\gamma(t)) \geq \alpha$.

Hence $c \geq \alpha$ and there exists $\bar{\varepsilon} > 0$ such that

$$c - \bar{\varepsilon} \geq \max \{ \varphi(u_0), \varphi(u_1) \}.$$

Assume that $K_c = \emptyset$. Let $\varepsilon \in]0, \bar{\varepsilon}[$ and $\gamma \in C([0,1], X, X)$ be given by lemma 1 where $U = \emptyset$. The definition of c implies the existence of $\gamma \in \Gamma$ such that $\max_{t \in [0,1]} \varphi(\gamma(t)) \leq c + \varepsilon$.

Let us now consider the continuous path $\tilde{\gamma}$ defined by $\tilde{\gamma}(t) = \gamma(1, \gamma(t))$. Since $\varphi(u_0) \leq c - \bar{\varepsilon}$, we have

$$\tilde{\gamma}(0) = \gamma(1, \gamma(0)) = \gamma(1, u_0) = u_0.$$

Similarly $\tilde{\gamma}(1) = u_1$. Thus $\tilde{\gamma} \in \Gamma$. But $\gamma([0,1]) \subset \varphi^{c+\varepsilon}$, so that $\gamma(1, \gamma([0,1])) \subset \varphi^{c-\varepsilon}$ and

$$c \leq \max_{t \in [0,1]} \varphi(\tilde{\gamma}(t)) \leq c - \varepsilon,$$

a contradiction. Hence $K_c \neq \emptyset$ and the proof is complete. $\square$

Let us now consider an other geometric situation. (Rabinowitz saddle point theorem).

Theorem 4. Assume that $g \in C^1(X, \mathbb{R})$ satisfying (PS) and that $X = Y \oplus Z$ with Y finite dimensional. If there exists constants $b_1 < b_2$ and a neighborhood Ω of 0 in Y bounded such that

$$g|_Z \geq b_2$$

$$g|_{\partial\Omega} \leq b_1$$

then g has a critical point.

Proof. Let us define

$$\Gamma = \{ \gamma \in C(\bar{\Omega}, X) : \gamma(y) = y \text{ if } y \in \partial\Omega \}.$$
$$c = \inf_{\gamma \in \Gamma} \max_{y \in \bar{\Omega}} g(\gamma(y)).$$

If P denotes the projector on Y with kernel Z , then $P\gamma \in C(\bar{\Omega}, Y)$ and $P\gamma(y) = y$ whenever $\gamma \in \Gamma$ and $y \in \partial\Omega$. By a topological degree argument, $P\gamma$ has a zero inside Ω , i.e. there exists $y \in \Omega$ such that $\gamma(y) \in Z$. (In the easy case when $\dim Y = 1$, the intermediate value theorem suffices.) In particular

$$\max_{y \in \Omega} g(\gamma(y)) \geq \inf_{z \in Z} g(z) \geq b_2.$$

Hence $c \geq b_2$ and there exists $\varepsilon > 0$ such that $c - \varepsilon \geq b_1$. Assume that $K_c = \emptyset$. Let $\varepsilon \in]0, \varepsilon[$ and $\gamma \in C([0, 1] \times X, X)$ be given by lemma 1 with $U = \emptyset$. Then exists $\gamma \in \Gamma$ such that

$$\max_{y \in \bar{\Omega}} g(\gamma(y)) \leq c + \varepsilon.$$

Set $\tilde{\gamma} = \gamma(1, \gamma(\cdot))$. If $y \in \partial\Omega$, then $g(y) \leq b_1$ and $\tilde{\gamma}(y) = \gamma(1, \gamma(y)) = \gamma(1, y) = y$. Thus $\tilde{\gamma} \in \Gamma$.

Since $\gamma(\bar{x}) \in \mathcal{Q}^{c+2}$, we have that
 $\gamma(1, \gamma(\bar{x})) \in \mathcal{Q}^{c-2}$ and
 $c \leq \max_{y \in \bar{x}} \mathcal{Q}(\tilde{\gamma}(y)) \leq c-2$,

a. contradiction. Hence $K_c \neq \emptyset$. $\square$

The following application is due essentially to Ahmad, Lazer and Paul.

Theorem 5. Assume that V and $D_x V$ are continuous on $[0, T] \times \mathbb{R}^N$ and that there exists a constant $a > 0$ such

$$|D_x V(t, x)| \leq a$$

on $[0, T] \times \mathbb{R}^N$. If

$$(10) \quad \int_0^T V(t, x) dt \rightarrow +\infty, \text{ as } |x| \rightarrow \infty$$

then problem (5) has at least one solution.

Proof. We shall apply theorem 4 to $\mathcal{Q}(u) = \int_0^T \left[\frac{1}{2} |\dot{u}(t)|^2 - V(t, u(t)) \right] dt$ with $X = H_T^1$,

$$Y = \{y: [0, T] \rightarrow \mathbb{R}^N : y \text{ is constant}\}$$

$$Z = \{z \in H_T^1 : \int_0^T z(t) dt = 0\}.$$

1) $\mathcal{Q}$ is bounded from below on Z .

$$\text{For } z \in Z \text{ we have, using Sobolev inequality,}$$

$$\mathcal{Q}(z) = \frac{\| \dot{z} \|^2}{2} - \int_0^T V(t, 0) dt - \int_0^T \int_0^1 (D_x V(t, s z(t)), z(t)) ds dt$$

$$\geq \frac{\| \dot{z} \|^2}{2} - c_1 - T a \| z \|_\infty$$

$$\geq \frac{\| \dot{z} \|^2}{2} - c_1 - c_2 \| \dot{z} \|_{L^2},$$

Thus $b_Z = \inf_Z \mathcal{Q} > -\infty$.

2) φ satisfies the (PS) condition.

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Let (u_k) be a sequence in H_T^1 such that $\varphi(u_k)$ is bounded and $\varphi'(u_k) \rightarrow 0$ as $k \rightarrow \infty$. As is known 2, it suffices to prove that (u_k) is bounded in H_T^1 .

There is some k_0 such that, for $k \geq k_0$,

$$|\langle \varphi'(u_k), h \rangle| \leq \|h\|, \quad \forall h \in H_T^1.$$

In particular, by Holder inequality,

$$c_1 \|\dot{u}_k\|_{L^2} \geq \|\tilde{u}_k\| \geq |\langle \varphi'(u_k), \tilde{u}_k \rangle| =$$

$$= \int_0^T [|\dot{u}_k(t)|^2 - \langle \nabla F(t, u_k(t), \tilde{u}_k(t)) \rangle] dt$$

$$\geq \|\dot{u}_k\|_{L^2}^2 - a_T \|\tilde{u}_k\|_\infty$$

$$\geq \|\dot{u}_k\|_{L^2}^2 - c_2 \|\dot{u}_k\|_{L^2}.$$

Thus $\dot{u}_k$ is bounded in L^2 and $\tilde{u}_k$ is bounded in H_T^1 . Finally

$$\begin{aligned} c_3 \leq \varphi(u_k) &= \int_0^T \left[\frac{1}{2} |\dot{u}_k(t)|^2 + V(t, \tilde{u}_k(t)) - V(t, u_k(t)) - V(t, \tilde{u}_k(t)) \right] dt \\ &\leq c_4 - \int_0^T V(t, \tilde{u}_k) dt \end{aligned}$$

and $\tilde{u}_k$ is bounded by assumption (10).

3) End of the proof.

If $y \in Y$ then $\varphi(y) = - \int_0^1 V(t, y) dt$. By assumption (10) there exists $R > 0$ such that

$$\|y\| \geq R \Rightarrow \varphi(y) \leq b_2 - 1 =: b_1.$$

It suffices then to apply theorem 4 to get a critical point of φ and, consequently, a solution of problem (1). $\square$

III - Lusternik-Schnirelman Theory.

A basic method in critical point theory is to obtain the existence of multiple critical points by using invariance properties of the functional. In this section, we consider the case of a "periodic" functional defined on a Banach space X .

We shall need the following definitions.

A subset C of a topological space Y is contractible in Y if there exists $h \in C([0,1], C, Y)$ and $y \in Y$ such that

$$h(0, u) = u, h(1, u) = y, \quad \forall u \in C.$$

A subset A of a topological space Y has category k in Y if k is the least integer such that A can be covered by k closed sets contractible in Y . The category of A in Y is denoted by $\text{cat}_Y(A)$.

A metric space Y is an absolute neighborhood extensor, shortly an ANE, if for every metric space E , every closed subset F of E and every $f \in C(F, Y)$ there exists a continuous extension of f defined on a neighborhood of F in E .

Examples. [45].

a) $\text{cat}_{S^N}(S^N) = 2, \quad \text{cat}_{T^N}(T^N) = N+1.$

b) A finite product of ANE is an ANE.

A convex subset of a normed space is an ANE.

A circle is an ANE.

Proposition 4. Let Y, Z be topological spaces. 21
and let $A, B \subset Y$.

- i) If $A \subset B$, then $\text{cat}_Y(A) \leq \text{cat}_Y(B)$ (monotonicity),
- ii) $\text{cat}_Y(A \cup B) \leq \text{cat}_Y(A) + \text{cat}_Y(B)$ (subadditivity),
- iii) If A is closed and $B = \gamma(1, A)$, where $\gamma \in C([0, 1], A, Y)$ is such that $\gamma(0, u) = u$ for every $u \in A$, then $\text{cat}_Y(A) \leq \text{cat}_Y(B)$.
- iv) If A is a closed subset of an ANE, Y then there exists a closed neighborhood U of A such that $\text{cat}_Y(A) = \text{cat}_Y(U)$ (continuity).

Proof. see [4]. $\square$

Let G be a discrete subgroup of a Banach space X and let $\pi: X \rightarrow X/G$ be the canonical projection. A subset A of X is G -invariant if $\pi^{-1}(\pi(A)) = A$. A function f defined on X is G -invariant if $f(u+g) = f(u)$ for every $u \in X$ and every $g \in G$. If a differentiable functional $\phi: X \rightarrow \mathbb{R}$ is G -invariant, then ϕ' is also G -invariant. Consequently, if u is a critical point of ϕ , then $\pi^{-1}(\pi(u))$ is a set of critical points of ϕ , and is called a critical orbit of ϕ .

A G -invariant differentiable function $\phi: X \rightarrow \mathbb{R}$ satisfies the (PS) c condition if, for every sequence (u_k) in X such that $\phi(u_k)$ is bounded and $\phi'(u_k) \rightarrow 0$, the sequence $(\pi(u_k))$ contains a convergent subsequence.

Theorem 5. Let $\varphi \in C^1(X, \mathbb{R})$ be a G -invariant functional satisfying the $(PS)_c$ condition. If φ is bounded from below and if the dimension N of the space generated by G is finite, then φ has at least $N+1$ critical orbits.

Lemma 2. Under the assumption of theorem 5, if U is an open invariant neighborhood of K_c then, for every $\varepsilon > 0$ there exists an $\varepsilon \in]0, \varepsilon[$ and $\gamma \in C([0, 1] \times X, X)$ satisfying properties (a), (b), (c) of lemma 1. Moreover (d) $\gamma(t, u+g) = \gamma(t, u) + g, \forall t \in [0, 1], \forall u \in X, \forall g \in G$.

Proof. We consider the special case of X Hilbert space and $\varphi \in C^2(X, \mathbb{R})$.

It is easy to construct the vector field f as in lemma 1. Moreover the invariance of U and $\nabla \varphi$ implies the invariance of f . A solution to the Cauchy problem

$$\dot{\sigma} = f(\sigma)$$

$$\sigma(0) = u + g \quad (g \in G)$$

is given by $\sigma(t, u) + g$. By uniqueness $\sigma(t, u+g) = \sigma(t, u) + g$. Since $\gamma(t, u) = \sigma(\sqrt{t}t, u)$ the proof is complete. $\square$

Lemma 3. Under the assumption of theorem 5, there exists a closed invariant subset Ω of X such that $\dim_{T(x)} \pi(\Omega) = N+1$.

Proof. Let V be the space generated by G . Then X is isomorphic to $\mathbb{R}^N \times Z$, where Z is a complement of V , G is isomorphic to $\mathbb{Z}^N$ and $\pi(x)$ is isomorphic to $T^N \times Z$.

Setting $A = \mathbb{R}^N \times \{0\}$, we obtain

$$\text{cat}_{\pi(x)}(\pi(A)) = \text{cat}_{T^N \times X}(T^N \times \{0\}) = \text{cat}_{T^N}(T^N) = N+1. \square$$

By lemma 3, for $1 \leq j \leq N+1$, the following set is non-empty:

$$\mathcal{A}_j = \{A \subset X : A \text{ is closed, invariant and } \text{cat}_{\pi(x)}(\pi(A)) \geq j\}.$$

Define

$$c_j = \inf_{A \in \mathcal{A}_j} \sup_A g.$$

Since $\mathcal{A}_j \subset \mathcal{A}_{j-1}$, we have

$$-\infty < \inf_X g \leq c_1 \leq c_2 \leq \dots \leq c_{N+1} < +\infty.$$

Lemma 4. Under the assumptions of theorem 6, if $c_k = c_j = c$ for some $1 \leq j \leq k \leq N+1$ then $\text{cat}_{\pi(x)}(\pi(K_c)) \geq k-j+1$.

Proof. Since $\pi(x) : T^N \times X \rightarrow \mathbb{R}^N \times X$ is an ANE and since $\pi(K_c)$ is compact by (PS)_c, then exists, by proposition 4, a closed neighborhood N of $\pi(K_c)$ such that $\text{cat}_{\pi(x)}(N) = \text{cat}_{\pi(x)}(\pi(K_c))$. Let $\varepsilon > 0$ be given by lemma 2 applied to $\tilde{\varepsilon} = 1$ and $U = \text{int } \pi^{-1}(N)$. By definition of c_k , there exists $A \in \mathcal{A}_k$ such that $\max_A g \leq c_k + \varepsilon = c + \varepsilon$.

Proposition 4 implies that

$$\begin{aligned} (11) \quad k &\leq \text{cat}_{\pi(x)}(\pi(A)) \leq \text{cat}_{\pi(x)}(\pi(A \setminus U) \cup \pi(U)) \\ &\leq \text{cat}_{\pi(x)}(\pi(A \setminus U)) + \text{cat}_{\pi(x)}(N) \\ &= \text{cat}_{\pi(x)}(\pi(A \setminus U)) + \text{cat}_{\pi(x)}(\pi(K_c)). \end{aligned}$$

By lemma 2, $c_3(1, A \setminus U) < c - \varepsilon$ and γ induces

a deformation from $\pi(R \setminus U)$ to $\pi(C)$.

But $\max_C \varphi \leq c - \varepsilon = c_j - \varepsilon$. By the definition

of c_j , $\text{cat}_{\pi(X)}(\pi(C)) \leq j-1$, so that, by proposition 4, $\text{cat}_{\pi(X)}(\pi(R \setminus U)) \leq j-1$.

Hence (11) implies that $k \leq j-1 + \text{cat}_{\pi(X)}(\pi(K_c))$. $\square$

Proof of theorem 6. If $c_j < c_k$ whenever $1 \leq j < k \leq N+1$, then φ has at least $N+1$ critical values.

If $c_j = c_k = c$ for some $1 \leq j < k \leq N+1$ then $\text{cat}_{\pi(X)}(\pi(K_c)) \geq 2$ and K_c

contains infinitely many critical orbits. $\square$

Theorem 6 generalizes a recent result of S. Li (ITCP Rep. IC-85-191) and was also motivated by a recent result due independently to J. Krasinkin and P. Rabinowitz which is the following:

Consider the classical Lagrangian $L(t, x, y) = \frac{1}{2} (M(t, x) y, y) - V(t, x) + f(t, x)$.

Assume that

(P1) $M \in C^1(\mathbb{R} \times \mathbb{R}^N, \mathcal{L}(\mathbb{R}^N))$ is a symmetric matrix which is T -periodic in t and x_i , $1 \leq i \leq N$, (for simplicity).

(P2) There exists $\alpha > 0$ such that $(M(t, x) y, y) \geq \alpha |y|^2$, $\forall [t, x, y] \in \mathbb{R} \times \mathbb{R}^N \times \mathbb{R}^N$.

(R3) $\forall \epsilon \in C^1(\mathbb{R} \times \mathbb{R}^N, \mathbb{R})$ is T -periodic in t and $\epsilon_i, 1 \leq i \leq N$.

(R4) $f \in C(\mathbb{R}, \mathbb{R}^N)$ is T -periodic in t and $\int_0^T f(t) dt = 0$.

Then the Lagrangian system

$$\frac{d}{dt} D_x L(t, u(t), \dot{u}(t)) = D_x L(t, u(t), \dot{u}(t))$$

has at least $N+1$ geometrically distinct periodic solutions.

In particular, when the mean value of the forcing term is 0, a simple forced pendulum has at least 2 geometrically distinct periodic solutions and a double forced pendulum has at least 3 geometrically distinct periodic solutions.

The above result follows easily from theorem 6. In particular, if (e_i) denotes the canonical basis of $\mathbb{R}^N$, the functional

$$\varphi(u) = \int_0^T L(t, u(t), \dot{u}(t)) dt$$

satisfies the following periodicity property on H_T^1 :

$$\varphi(u + Te_i) = \varphi(u), \quad 1 \leq i \leq N.$$

Remark. Using Ekeland principle it is possible to extend Lusternik-Schnirelman theory to C^1 Finsler-manifolds (Szulkin, 1987).

IV - Bourbaki-Ulam theorem and index theory.

In this section we consider functionals which are invariant under the action of a compact topological group.

Let G be a topological group. A representation of G over a Banach space X is a family $\{T(g)\}_{g \in G}$ of linear operators $T(g): X \rightarrow X$ such that

$$T(1) = \text{id}$$

$$T(g_1 + g_2) = T(g_1) \circ T(g_2)$$

$(g, u) \mapsto T(g)u$ is continuous.

A subset A of X is invariant (under the representation) if $T(g)A = A$ for all $g \in G$.

A representation is isometric if $\|T(g)u\| = \|u\|$ for all $g \in G$ and all $u \in X$.

Examples. 1) $G = \mathbb{Z}_2$, $T(1) = \text{id}$, $T(2) = -\text{id}$.

2) $G = S^1 \simeq \mathbb{R}/\mathbb{T}$, $X = W^{1,2}$, $T(\theta)$ is given by $(T(\theta)u)(t) = u(t + \theta)$.

Let G be a compact topological group and let $\{T(g)\}_{g \in G}$ be an isometric representation of G over a Banach space X .

A mapping R between two invariant subsets of X is equivariant if

$$R \circ T(g) = T(g) \circ R, \quad \forall g \in G.$$

A real function ϕ defined on an invariant subset of X is invariant if

$$\phi \circ T(g) = \phi, \quad \forall g \in G.$$

An index (for $\{T(g)\}_{g \in G}$) is a mapping from the closed invariant subsets of X into $\mathbb{N} \cup \{\infty\}$ such that

$$(i) \text{ ind } A = 0 \Leftrightarrow A = \emptyset$$

(ii) if $R: A_1 \rightarrow A_2$ is equivariant and continuous, then $\text{ind } A_1 \leq \text{ind } A_2$.

(iii) if A is compact and invariant, then exists a closed invariant neighborhood N of A such that $\text{ind } N = \text{ind } A$.

(iv) $\text{ind } (A_1 \cup A_2) \leq \text{ind } A_1 + \text{ind } A_2$ for invariant subsets A_1 and A_2 .

Examples.

1) (Krasnoselski) let $G = \mathbb{Z}_2$, $T(0) = \text{id}$, $T(1) = -\text{id}$.

The $\mathbb{Z}_2$ -index of a closed invariant subset A of X is the smallest integer k such that there exists an odd mapping $\Phi \in C(A, \mathbb{R}^k \setminus \{0\})$. If such a mapping does not exist, we define $\text{ind } A = +\infty$. Finally $\text{ind } \emptyset = 0$.

2) (Benci) . let $\{T(\theta)\}_{\theta \in S^1}$ be an isometric representation of S^1 over X .

The S^1 -index of a closed invariant subset A of X is the smallest integer k such that there exists a $n \in \mathbb{N} \setminus \{0\}$ and a $\Phi \in C(A, \mathbb{C}^k \setminus \{0\})$ satisfying the following equivariance property:

$$\Phi(T(\theta)u) = e^{in\theta} \Phi(u), \quad \theta \in S^1, u \in A.$$

If such a mapping does not exist, we define $\text{ind } A = +\infty$. Finally $\text{ind } \emptyset = 0$.

The verification of properties (i) to (iv) is left to the reader (see [47, 57]).

The following deep results are formulation of Borsuk-Ulam theorems.

Let us recall that

$$\text{Fix}(G) = \{u \in X : T(g)u = u, \forall u \in X\}.$$

Theorem 7. The $\mathbb{Z}_2$ -index of the sphere S^{k-1} is k .

Theorem 8. Let $|T(\theta)|_{0 \leq s \leq 2}$ be an isometric representation of S^1 on $\mathbb{R}^{2k}$ such that $\text{Fix}(\theta) = \{0\}$. Then the S^1 -index of the sphere S^{2k-2} is k .

A proof using transversality and degree theory is contained in [4].

Let us now consider the nonlinear eigenvalue problem

$$(12) \quad \varphi'(u) = \mu \chi'(u), \quad \chi(u) = a,$$

where $\varphi, \chi \in C^1(\Omega, \mathbb{R})$, Ω open subset of $\mathbb{R}^N$. Assume that

$$\chi(u) = a \Rightarrow \chi'(u) \neq 0$$

so that

$$Z_a = \{u \in \Omega : \chi(u) = a\}$$

is a submanifold of $\mathbb{R}^N$ of codimension 1. Equation (12) is equivalent to

$$\text{Ker } \chi'(u) \subset \text{Ker } \varphi'(u).$$

Since the tangent space to Z_a at u , $T_u Z_a$ is given by $\text{Ker } \chi'(u)$, equation (12) is equivalent to

$$T_u Z_a \subset \text{Ker } \varphi'(u)$$

so that u is a solution of (12) if and only if u is a critical point of φ restricted to Z_a .

By the Lagrange multiplier rule, a maximum or a minimum of $\varphi|_{Z_a}$ is a critical point of $\varphi|_{Z_a}$.

Assume now that Ω and φ are invariant under an isometric representation of the compact group G over $\mathbb{R}^N$. It is then easy to verify that $\nabla \varphi$ is equivariant, i.e.

$$\nabla \varphi \circ T(g) = T(g) \circ \nabla \varphi, \quad \forall g \in G$$

Consider the nonlinear eigenvalue problem

$$\varphi'(u) = \mu u \quad |u|^2 = 1$$

If u is a critical point of $\varphi|_{S^{N-1}}$ then $\{T(g)u : g \in G\}$ is a set of critical points of $\varphi|_{S^{N-1}}$ and is called a critical orbit of $\varphi|_{S^{N-1}}$.

Lusternik and Schnirelman showed that if φ is $\mathbb{Z}_2$ -invariant (i.e. even), then $\varphi|_{S^{N-1}}$ has at least N distinct pairs

$\{u, -u\}$ of critical points. This generalizes the classical result on the existence of N linearly independent eigenvectors for every linear symmetric mapping $L: \mathbb{R}^N \rightarrow \mathbb{R}^N$.

In the case of S^1 -invariant functionals, we have the following result due to Krasnoselski.

Theorem 3. Let $T(\partial)$ be an isometric representation of S^1 over $\mathbb{R}^{2k}$ such that $\text{Fix}(S^1) = \{0\}$. Let $\Omega \subset \mathbb{R}^N$ be open and invariant, let $\phi \in C^1(\Omega, \mathbb{R})$ be invariant. If $S^{2k-1} \subset \Omega$ then $\phi|_{S^{2k-1}}$ has at least k critical orbits.

Proof. 1) let us define

$$\mathcal{O}_j = \{A \subset S^{2k-1} : A \text{ is closed, invariant and ind } A \geq j\}$$

By theorem 8, $\mathcal{O}_j \neq \emptyset$ for $1 \leq j \leq k$. Let

$$c_j = \inf_{A \in \mathcal{O}_j} \max_A \phi$$

so that $\min \phi|_{S^{2k-1}} \leq c_1 \leq c_2 \leq \dots \leq c_k \leq \max \phi|_{S^{2k-1}}$.

If

$$K_c = \{u \in S^{2k-1} : u \text{ is a critical point of } \phi|_{S^{2k-1}}\}$$

then it suffices to prove that

$$c_i = c_j = c, 1 \leq i < j \leq k \Rightarrow \text{ind } K_c \geq i - j + 1,$$

since the index of a finite number of S^1 -orbits is 1.

2) Consider the vector field defined on S^{2k-1} by

$$v(u) = -\nabla \phi(u) - (\nabla \phi(u), u) u.$$

Clearly v is equivariant and

$$(u, v(u)) = 0$$

$$(\nabla \phi(u), v(u)) = -(\nabla \phi(u), u) = -|v(u)|^2 \leq 0,$$

for every $u \in S^{2k-1}$. Assuming that

$\phi \in C^2(\Omega, \mathbb{R})$, it is then easy to

construct, as in lemma 1, a deformation

$\gamma \in C([0, 1] \times S^{2k-1}, S^{2k-1})$ satisfying the usual properties and such that

$$\gamma(t, T(\partial)u) = T(\partial)\gamma(t, u)$$

for every $t \in [0, 1]$, $\partial \in S^1$, $u \in S^{2k-1}$.

The end of the proof is similar to the proof of lemma 4 and is left to the reader. $\square$

Remarks. 1. In the C^1 case, it suffices to replace ϕ by an appropriate pseudo-gradient vector field.

2. Assume that $\chi \in C^1(\Omega, \mathbb{R})$ is invariant and that

$$Z_a = \{u \in \Omega : \chi(u) = a\}$$

is diffeomorphic to S^{2k-2} . If the corresponding diffeomorphism is equivariant, then theorem 3 implies the existence of k critical orbits for $\phi|_{Z_a}$.

3. Since S^{2k-2}/S^1 is not, in general, a manifold, it is impossible to deduce theorem 3 from Weinstein-Schwarzman theory.

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V. Bifurcation theory and periodic solutions of Hamiltonian systems near an equilibrium.

Consider the nonlinear eigenvalue problem
 (13) $\nabla \alpha(u) + \lambda \nabla \beta(u) = 0$
 where α and β are C^2 functionals defined on a Hilbert space X . We shall assume that

$$\nabla \alpha(0) = \nabla \beta(0) = 0.$$

Thus $\mathbb{R} \times \{0\}$ is a branch of trivial solutions of (13). Using Liapunov-Schmidt method and an elementary variational argument, we shall construct two distinct one parameter families of non trivial solutions of (13). We shall also prove a stronger multiplicity result when α and β are S^1 invariant.

Our basic assumptions are

(H1) $\nabla \alpha(0) = 0, \nabla \beta(0) = 0$

(H2) $L = \alpha''(0)$ is Fredholm, i.e. $\dim \text{Ker } L < +\infty$, $\text{codim } \text{R}(L) < +\infty$.

(H3) $\dim \text{Ker } L \geq 2$ and $M = \beta''(0)$ is positive definite on $\text{Ker } L$.

Remarks. 1. Since L is Fredholm, $\text{R}(L)$ is closed. The symmetry of L implies that X is the orthogonal sum of $\text{R}(L)$ and $\text{Ker } L$.

2. Without loss of generality, we can assume that $\alpha(0) = \beta(0) = 0$.

3. Let P (resp. Q) be the orthogonal projector on $\text{Ker } L$ (resp. $\text{R}(L)$). Clearly $Q = 1 - P$.

We shall use the following notations:

$$A = \nabla \alpha, B = \nabla \beta, R = A - L, S = B - M.$$

Theorem 10. (C. Stuart). Under the assumptions $(H_1, 2, 3)$, for every $\varepsilon > 0$ small enough, there exists at least two solutions $[\lambda(\varepsilon), u(\varepsilon)]$ of (13) such that $\beta(u(\varepsilon)) = \varepsilon$. Moreover $\lambda(\varepsilon) \rightarrow 0$, $\varepsilon \rightarrow 0$.

Proof. 1) Liapunov-Schmidt reduction.

Equation (13) is equivalent to the system

$$P[R(u+w) + \lambda B(u+w)] = 0$$

$$(14) \quad Lw + Q[R(u+w) + \lambda B(u+w)] = 0,$$

where $u = Pu$, $w = Qu$.

Since $L: R(L) \rightarrow R(L)$ is invertible, it follows from the implicit function theorem that equation (14) defines near $\lambda=0$, $u=0$, $w=0$ a C^1 function $w = w^*(\lambda, 0)$. Since $\mathbb{R} \times \{0\}$ is a branch of solutions, $w^*(\lambda, 0) = 0$. Differentiating $Lw^*(\lambda, 0) + Q[R(u+w^*, \lambda, 0) + \lambda B(u+w^*, \lambda, 0)] = 0$ with respect to λ at $[0, 0]$ and using the fact that $R'(0) = 0$, we obtain

$$L D_\lambda w^*(0, 0) = 0,$$

i.e. $D_\lambda w^*(0, 0) = 0$. Thus

$$\|w^*(\lambda, 0)\| / \|\lambda\| \rightarrow 0, \quad \|\lambda\| \rightarrow 0$$

uniformly for λ near 0.

Near $\lambda=0$, $u=0$, $w=0$, equation (13) is now equivalent to

$$(14) \quad P[R(u+w^*(\lambda, u)) + \lambda B(u+w^*(\lambda, u))] = 0.$$

2) Deparametrization. Taking the inner

product of (14) with u , we obtain

$$(15) \quad (R(u+w^*(\lambda, u)), u) + \lambda (B(u+w^*(\lambda, u)), u) = 0.$$

It follows from a careful application of the implicit function theorem that equation (15) defines, near 0, for $u \neq 0$, a C^1 function

$\lambda = \lambda^*(v)$. Moreover λ^* is extended continuously at 0 by setting $\lambda^*(0) = 0$.

Let us define, near 0,

$$f(v) = v + (\lambda^*(v), v)$$

so that

$$\|f(v)\|/\|v\| \rightarrow 0, \quad \|v\| \rightarrow 0.$$

Equation (13) is now equivalent to

$$(16) \quad P[R(v + f(v)) + \lambda^*(v) B(v + f(v))] = 0.$$

3) Constrained extremization. For $\rho > 0$ small enough, the function

$$X(v) = B(v + f(v))$$

is continuous on $B(0, \rho)$ and C^1 on $B(0, \rho) \setminus \{0\}$.

It follows from (13) and some careful estimates (see [4]), that for $\rho > 0$ small enough

$$(17) \quad \begin{aligned} X(v) &\geq c\|v\|^2, & \forall v \in B(0, \rho) \\ (\nabla X(v), v) &\geq c\|v\|^2, & \forall v \in B(0, \rho) \setminus \{0\} \end{aligned}$$

where $c > 0$ is a constant. Hence, for $\varepsilon > 0$ small enough,

$$Z_\varepsilon = \{v \in \text{Ker } L : X(v) = \varepsilon\}$$

is a compact subset of $B(0, \rho) \setminus \{0\}$ and

$$(18) \quad v \in Z_\varepsilon \Rightarrow (\nabla X(v), v) > 0.$$

The function

$$\varphi(v) = \alpha(v + f(v)) + \lambda^*(v)(X(v) - \varepsilon)$$

restricted to Z_ε achieves his minimum

at a point v_ε . Since $\nabla X(v_\varepsilon) \neq 0$, the Lagrange multiplier rule implies the existence of μ such that

$$\nabla \varphi(v_\varepsilon) = \mu \nabla X(v_\varepsilon),$$

i.e.

$$(R(v_\varepsilon + f(v_\varepsilon)) + \lambda^*(v_\varepsilon) B(v_\varepsilon + f(v_\varepsilon)), h + f'(v_\varepsilon)h) = \mu (\nabla X(v_\varepsilon), h),$$

for every $h \in \text{Ker } L$. Since $f'(v_\varepsilon)h \in R(\varphi)$, the

definition of $f(u_2) = w^*(\lambda^*(u_2), u_2)$ implies that

$$(R(u_2 + f(u_2)) + \lambda^*(u_2) B(u_2 + f(u_2)), f(u_2)h) = 0.$$

Using the fact that $PA = PR$, we have, for $h \in \text{Ker } L$,

$$(R(u_2 + f(u_2)) + \lambda^*(u_2) B(u_2 + f(u_2)), h) = \mu(\nabla \chi(u_2), h).$$

It follows from the definition of λ^* that

$$0 = (R(u_2 + f(u_2)) + \lambda^*(u_2) B(u_2 + f(u_2)), u_2) = \mu(\nabla \chi(u_2), u_2).$$

We obtain from (18) that $\mu = 0$. Finally

$$(R(u_2 + f(u_2)) + \lambda^*(u_2) B(u_2 + f(u_2)), h) = 0$$

for every $h \in \text{Ker } L$, i.e. u_2 is a solution of (16). Thus

$$[\lambda(z) = \lambda^*(u_2), u_2 = u_2 + f(u_2)]$$

is a solution of (13) such that

$$B(u_2) = \chi(u_2) = \varepsilon.$$

By (17), $u_2 \rightarrow 0$ as $\varepsilon \rightarrow 0$, so that $\lambda(\varepsilon) = \lambda^*(u_2) \rightarrow \lambda^*(0) = 0$ as $\varepsilon \rightarrow 0$. The other solution is obtained by maximizing φ on Σ_ε . $\square$

Let now $T(S)$ be an isometric representation of S^1 over the Hilbert space X such that

(H4) α and β are S^1 -invariant and $\text{Ker } L \cap \text{Fix}(S^1) = \{0\}$.

Remark. By (H4), A, B, L, M, R, S, P and Q are equivariant, $R(L)$ and $\text{Ker } L$ are invariant. Since $\text{Ker } L \cap \text{Fix}(S^1) = \{0\}$, the dimension of $\text{Ker } L$ is even.

Theorem 11. Under the assumptions of theorem 10 and (H4), for every $\varepsilon > 0$ small enough, there exists at least $\frac{1}{2} \dim \text{Ker } S^1$ solutions of (13) such that $\|u_\varepsilon\|_{C^1} \rightarrow 0$ as $\varepsilon \rightarrow 0$.

Proof. It is easy to verify that the invariance of $\mathcal{Q}$ and β implies the invariance of $\mathcal{P}$ and χ . By (17), for $\varepsilon > 0$ small enough,

$$Z_\varepsilon = \{v \in \text{Ker } L : \chi(v) \leq \varepsilon\}$$

is diffeomorphic to S^{2k-1} when $2k = \dim \text{Ker } L$. Moreover the corresponding diffeomorphism is equivariant. Hence theorem 3 implies the existence of k critical orbits for $\mathcal{P}|_{Z_\varepsilon}$.

It is then easy to complete the proof as in theorem 10. $\square$

Let $H \in C^2(\mathbb{R}^{2N}, \mathbb{R})$ be such that $H(0) = 0$, $\nabla H(0) = 0$. We consider the existence of periodic orbits of the Hamiltonian system

$$(18) \quad J\dot{u}(t) + \nabla H(u(t)) = 0$$

on the energy surface $H^{-1}(\varepsilon)$ for small ε .

After the classical results of Liapunov (1900), depending on a non resonance condition, a significant progress was made by Weinstein around 1970. We shall deduce Weinstein theorem from a generalization due to Moser.

Theorem 12. (Moser). Assume that the linearized system

$$(20) \quad J\ddot{u}(t) + H''(0) u(t) = 0$$

has $2k$ linearly independent T -periodic solutions and that, if $u \neq 0$ is a T -periodic solution of (20), then

$$(H''(0) u(t), u(t)) = (H''(0) u(0), u(0)) > 0.$$

Then, for every $\varepsilon > 0$ small enough, there exists at least k periodic orbits of (19) on $H^{-1}(\varepsilon)$ whose period are near T .

Remarks. 1. By assumption $H''(0)$ is non singular.

2. The period T is not necessarily the minimal period of the corresponding solutions of (19) and (20).

Proof. We normalize the problem by fixing the period at 1. After the change of variable $s = T^{-1}t$, (19) becomes (21)

$$J\dot{z}(s) + T \nabla H(z(s)) = 0.$$

Any 1-periodic solution of (21) corresponds to a T -periodic solution of (19). Setting $\tau = T + \lambda$ we obtain the bifurcation problem

$$\nabla \alpha(z) + \lambda \nabla \beta(z) = 0$$

where the functionals

$$\alpha(z) = \int_0^1 [T \dot{z}(s), z(s)] ds + TH(z(s)) ds$$

$$\beta(z) = \int_0^1 H(z(s)) ds,$$

are defined on the Hilbert space $X = W^{1,2}_1 = H^1_1$. The functionals α and β are C^2 and invariant

representation of $S^1 = \mathbb{R}/\mathbb{Z}$ defined on X by the translations in time.

Let $L = \alpha''(0)$. Then $z \in \text{Ker } L$ if and only if

$$J \dot{z}(1) + T H''(0) z(1) = 0$$

$$z(0) = z(1)$$

By assumption, $\dim \text{Ker } L = 2k$ and $H = \beta''(0)$ is positive definite on $\text{Ker } L$. Since $H''(0)$ is non singular, $\text{Ker } L \cap \text{Fix}(S^1) = \{0\}$. It suffices then to apply theorem 11. $\square$

Corollary (Weinstein). Assume $H \in C^2(\mathbb{R}^{2N}, \mathbb{R})$, with $H(0) = 0$, $\nabla H(0) = 0$ and $H''(0)$ positive definite. For every $\varepsilon > 0$ small enough, there exists at least N period orbits of (19) on $H^{-1}(\varepsilon)$.

Proof. By assumption (20) has $2N$ linearly independent periodic solutions. The periodic solutions of (20) split into n families with incommensurable periods $T_1, \dots, T_n$ and dimensions $k_1, \dots, k_n$. We have that $k_1 + \dots + k_n = 2N$. Theorem 12, applied to each of these families, implies the existence of $k_1/2 + \dots + k_n/2 = N$ periodic orbits of $H^{-1}(\varepsilon)$ for $\varepsilon > 0$ small enough. These periodic orbits are distinct within the same family and from one family to another, because they have no common period for $\varepsilon > 0$ small enough. $\square$

Examples.

1) Linear case. Assume that the α_i 's are positive and incommensurable. Then each energy surface of the quadratic Hamiltonian

$$H(u) = \sum_{i=1}^N \alpha_i (u_i^2 + u_{i+N}^2)$$

contains exactly N periodic orbits.

2) Heron - Heiles Hamiltonian. The following Hamiltonian appears in astrophysics:

$$H(u) = (u_1^2 + u_2^2 + u_3^2 + u_4^2)/2 + u_1^2 u_2 - u_2^3/3.$$

For $\varepsilon > 0$ small enough, $H^{-1}(\varepsilon)$ contains at least 2 periodic orbits.

3) Moser example. The only periodic solution corresponding to the following Hamiltonian is the equilibrium:

$$H(u) = (u_1^2 + u_3^2 - u_2^2 - u_4^2) + (u_1^2 + u_2^2 + u_3^2 + u_4^2)(u_3 u_4 - u_1 u_2)$$

(It suffices to differentiate $u_1 u_4 + u_2 u_3$.) But all the orbits of the linearization at 0 are periodic.

Remark. See [27], [42], [50] for other aspects of variational methods in bifurcation theory.

Appendix.

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Proposition. Let X be a Banach space and let $q \in C^1(X, \mathbb{R})$. If q satisfies (P.S.) and is bounded from below, then q is coercive, i.e. $q(u) \rightarrow +\infty$ as $\|u\| \rightarrow \infty$.

Proof. If q is not coercive, then $c = \sup \{d \in \mathbb{R} : q^d \text{ is bounded}\}$ is a real number. Let U be an open bounded neighborhood of the compact set K_c and let $\varepsilon \in (0, 1)$ and γ be given by the deformation lemma. It follows from the definition of c that $q^{c+\varepsilon} \setminus U$ is unbounded and that $q^{c-\varepsilon}$ is bounded. On the other hand $\gamma(1, q^{c+\varepsilon} \setminus U) \subset q^{c-\varepsilon}$ and $\gamma(1, \cdot) : X \rightarrow X$ maps unbounded sets into unbounded sets. This is a contradiction. $\square$

Remark. The converse is valid in finite dimensions only.

I thank Mr. C. Amann for suggesting me this question.

