



Symplectic geometry and Floer homology

François Laudenbach

In memory of Laurent Schwartz whose charm was so very influential. I owe it to him to have become a mathematician.

Summary. Here are the revised notes for lectures held at the 13th Brazilian Topology Meeting in Belo Horizonte (July 2002). The purpose is to give an introduction to symplectic Floer homology and, in a simple case, a sketch of proof of the Arnold conjecture. This conjecture gives a lower bound for the number of fixed points of a Hamiltonian diffeomorphism in terms of the sum of the Betti numbers. Floer theory is a sort of infinite dimensional Morse theory on a loop space. The Morse index is replaced by the Maslov-Conley-Zehnder index. Some results about the Maslov cycle in the linear symplectic group are gathered in an appendix.

1 A conjecture by V.I. Arnold

Conjecture [4]. Actually the conjecture goes back to Arnold's note [1]. Here is the statement:

Let (M^{2n}, ω) be a symplectic closed manifold. Let φ_1 be a Hamiltonian diffeomorphism of M . Then φ_1 has at least as many fixed points as a smooth function has critical points:

$$\# \text{Fix } \varphi_1 \geq \inf_{f \in C^\infty(M, \mathbb{R})} \# \text{crit } f.$$

If we think only of the generic case, it is natural to look at Morse functions. Remembering the Morse inequalities, one can extract a weaker form of the conjecture.

Weak conjecture. *For a generic Hamiltonian diffeomorphism, the number of fixed points is at least the sum of the Betti numbers:*

$$\#\text{Fix } \varphi_1 \geq \sum_i \text{rank } H_i(M, \mathbb{Z}).$$

Conley and Zehnder proved the strong conjecture for $M = T^{2n}$ in [14]. A particular case of the weak Arnold conjecture – when (M, ω) is called *monotone* – has been proven by Andreas Floer [24]. Since that time other cases of the weak conjecture have been proven; the most general case has been solved by K. Fukaya and K. Ono [25], and also by G. Liu and G. Tian [35]. But as stated above the Arnold conjecture is still an open question, in particular in the non-generic case (except for T^{2n}). In this mini-course we would like to explain a very simple case of the weak conjecture where Floer’s ideas apply. First we recall a little from symplectic geometry as well as the prehistory of the subject which begins with Poincaré’s *last geometric theorem*.

1.1 Basic definitions in symplectic geometry

A *symplectic structure* on a manifold M is given by a 2-form $\omega \in \Omega^2(M)$, which satisfies the two following conditions:

- 1) $d\omega = 0$,
- 2) for every $x \in M$, $\omega(x) : T_x M \times T_x M \rightarrow \mathbb{R}$ is bilinear, skew-symmetric and nondegenerate.

The determinant of a skew-symmetric $n \times n$ matrix vanishes when n is odd. Then condition 2 forces the dimension of M to be even, $\dim M = 2n$, and the non-degeneracy condition just says that $\omega^n := \omega \wedge \dots \wedge \omega$ is a volume form on M . In particular M is oriented.

Example. $M = \mathbb{R}^{2n}$ with coordinates (p_i, q_i) , $\omega_0 = \sum dp_i \wedge dq_i$ and $\omega_0^n = n!(dp_1 \wedge dq_1 \wedge \dots \wedge dp_n \wedge dq_n)$.

Darboux’s Theorem. *Any $2n$ -dimensional symplectic manifold is locally isomorphic to $(\mathbb{R}^{2n}, \omega_0)$.*

Definition 1.1. *Thinking of \mathbb{R}^{2n} as the cotangent space $T^*\mathbb{R}^n = \mathbb{R}^n \times \mathbb{R}^n$, where the q_i ’s are the coordinates on the basis and the p_i ’s the coordinates on the fibre, the standard Liouville form is $\lambda_0 = \sum p_i dq_i$ and $d\lambda_0 = \omega_0$.*

Theorem 1.2. *The Liouville form λ_0 is “natural” in the following sense. Let $\varphi : U \rightarrow V$ be a (local) diffeomorphism of \mathbb{R}^n , $Q = \varphi(q)$. Let $\Phi : U \times \mathbb{R}^{n*} \rightarrow V \times \mathbb{R}^{n*}$ be its cotangent extension:*

$$\Phi(q, p) := \left(\varphi(q), {}^t(D\varphi(q))^{-1} \right).$$

Then $\Phi^(\sum P_i dq_i) = \sum p_i dq_i = {}^t p \cdot dq$.*

Proof. $\Phi^*({}^tP \cdot dQ) = {}^t p \cdot D\varphi(q)^{-1} \cdot D\varphi(q) \cdot dq$. \square

Corollary 1.3. *Let $M = T^*N$ be the cotangent space of N . There exists a 1-form, called the canonical Liouville form λ , which is given in coordinates as $p dq$. Its differential $\omega = d\lambda$ is closed and is given locally by $dp \wedge dq$. So ω is a symplectic form.*

Other examples: T^{2n} , $P^n(\mathbb{C})$, the orientable surfaces equipped with an area form.

Definition 1.4. *With a smooth function $H : \mathbb{R} \times \mathbb{R}^{2n} \rightarrow \mathbb{R}$, $(t, p, q) \mapsto H(t, p, q) = H_t(p, q)$, called a Hamiltonian (function), one can associate a Hamiltonian differential system:*

$$\begin{cases} \dot{p} = -\frac{\partial H_t}{\partial q} \\ \dot{q} = +\frac{\partial H_t}{\partial p} \end{cases}$$

Example. With the Hamiltonian $H = \frac{p^2}{2} + V(q)$ the following system is associated:

$$\begin{cases} \dot{p} = -\nabla V \\ \dot{q} = p \end{cases}$$

When H is time independent, $H(p(t), q(t))$ is constant along any solution of the system.

Notation 1.5. Let us identify \mathbb{R}^{2n} and $\mathbb{C}^n = \{p_1 + iq_1, \dots\}$. If J denotes the operator of multiplication by i in each complex factor, the Hamiltonian system from definition 1.4 reads:

$$\dot{x} = J \nabla H_t(x).$$

Definition 1.6. *On a symplectic manifold (M, ω) let $H : \mathbb{R} \times M \rightarrow \mathbb{R}$ be a smooth Hamiltonian function. We denote X_{H_t} the Hamiltonian vector field ω -dual of dH_t :*

$$i(X_{H_t})\omega = -dH_t.$$

For each t , X_t is also called the symplectic gradient of the function H_t and reads $X_t = \nabla_\omega H_t$.

The sign $-$ is due to the usual convention in classical mechanics. When $(M, \omega) = (\mathbb{R}^{2n}, \omega_0)$, then $\nabla_\omega = J \nabla$.

Proposition 1.7. *A time dependent vector field X_t is Hamiltonian iff it satisfies the two following conditions:*

- 1) $\varphi_t^* \omega = \omega$ where φ_t denotes its flow,
- 2) for any closed curve c , $\int_c i(X_t)\omega = 0$.

Proof. \Rightarrow 1) $\frac{d}{dt}\varphi_t^*\omega = \varphi_t^*(L_{X_t}\omega) = \varphi_t^*(i(X_t)d\omega + di(X_t)\omega) = \varphi_t^*(-ddH_t) = 0$.

2) $\int_C dH_t = 0$.

\Leftarrow By the above calculation the first condition implies that the 1-form $i(X_t)d\omega$ is closed. According to the second condition it is exact. \square

Definition 1.8. The time 1 map φ_1 from a Hamiltonian flow is said a Hamiltonian diffeomorphism. $\text{Ham}(M)$ is the set of Hamiltonian diffeomorphisms.

Exercise. $\text{Ham}(M)$ is a group.

Lemma 1.9. The same φ_1 can be obtained from a Hamiltonian function which is 1-periodic in time.

Proof. Let H_t be a Hamiltonian yielding the flow φ_t . Then we have

$$\frac{d}{dt}\varphi_{\alpha(t)} = \dot{\varphi}_{\alpha(t)}\alpha'(t) = X_{H_t}(\varphi_{\alpha(t)})\alpha'(t) = X_{\alpha'(t)H_t}(\varphi_{\alpha(t)}) .$$

Then $K_t := \alpha'(t)H_t$ is a Hamiltonian for the isotopy $\{\varphi_{\alpha(t)}\}$. If $\alpha : [0, 1] \rightarrow [0, 1]$ we have $\varphi_{\alpha(1)} = \varphi_1$. Moreover if $\alpha(t) = i$, $i = 0, 1$, for t close to i , $K_t = 0$ near 0 and 1 and then it can be extended periodically. \square

Note. In the rest of this text any Hamiltonian H under consideration will be 1-periodic in time, ie: $H : S^1 \times M \rightarrow \mathbb{R}$, where $S^1 = \mathbb{R}/\mathbb{Z}$. If φ_t denotes its flow, we have the following 1-1 correspondence:

$$\text{Fix } \varphi_1 \iff \{1\text{-periodic orbits of } X_{H_t}\} .$$

1.2 The Poincaré – Birkhoff Theorem

Theorem 1.10. Let f be an area preserving homeomorphism of the annulus $A := S^1 \times [-1, +1]$ which satisfies the twist condition. Then f has at least 2 fixed points. [7]

Moreover there are infinitely many periodic points, all geometrically distinct [12].

The universal cover is $\tilde{A} = \mathbb{R} \times [-1, +1]$. Let (p, q) be the coordinates in this product. On \tilde{A} , p becomes an angular coordinate. The area form is $\omega = dp \wedge dq$. A lift \tilde{f} to \tilde{A} reads

$$\begin{aligned} (p, q) &\mapsto (p', q') \\ p' &= p + a(p, q) \\ q' &= q + b(p, q). \end{aligned}$$

where a and b are two functions 1-periodic in p , such that b vanishes for $q = \pm 1$:

Definition 1.11. *The homeomorphism f satisfies the twist condition if there is some lift \tilde{f} to \tilde{A} with $a(p, +1) > 0$ and $a(p, -1) < 0$.*

We will give a proof in two particular cases. The first one is very particular, but it shows clearly how the fixed points problem is related to the critical points problem of functions. The second one could be called *the Poincaré case*. The idea of *generating functions* comes to play there.

Proof.

1) *The very particular case.* Here $X = \nabla_{\omega} h$ where $h : A \rightarrow \mathbb{R}$ is a function which is independent of time and constant along each component of the boundary, with a (usual) gradient pointing inside the annulus. If \vec{n} is an outgoing vector at each point of ∂A , we have $\omega(X, \vec{n}) > 0$ and then, if f is the time 1 map of the flow of X the twist condition is satisfied.

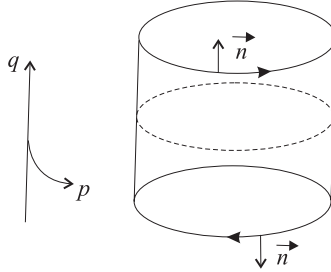


Figure 1

In this situation we have:

$$\text{crit } h \subset \text{Fix } f.$$

We have to prove: $\#\text{crit } h \geq 2$. Certainly there is some maximum which is not on the boundary since the gradient of h points in. So it is a first critical point lying in the interior of A . Let us assume there is only one maximum. The second critical point will be obtained by a minimax argument.

Let $\mathcal{T} = \{\alpha : [-1, +1] \rightarrow S^1 \times [-1, +1] \mid \alpha \text{ is } C^0, \alpha(\pm 1) \in S^1 \times \{\pm 1\}\}$. By definition the *minimax* value is

$$c := \inf_{\alpha \in \mathcal{T}} \left(\max_t h \circ \alpha(t) \right).$$

Due to the boundary condition and because the absolute maximum is unique, we have:

$$\max(h(S^1 \times \{\pm 1\})) < c < \max h.$$

The existence of the second critical point follows by the following lemma.

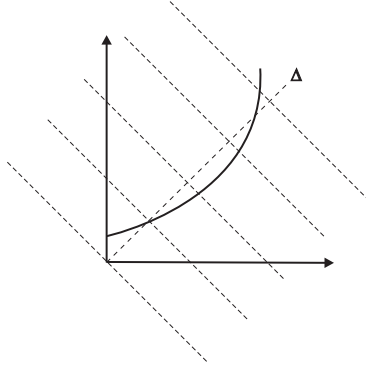


Figure 2

Lemma 1.12. *The minimax value c is a critical value.*

Proof. If not, for some $\varepsilon > 0$, $h^{-1}([c - \varepsilon, c + \varepsilon])$ is diffeomorphic to $h^{-1}(c) \times [c - \varepsilon, c + \varepsilon]$ and any path from \mathcal{T} lying in $h^{-1}([-\infty, c + \varepsilon])$ can be pushed by a homotopy with fixed endpoints into $h^{-1}([-\infty, c - \varepsilon])$ which contradicts the definition of c . □

2) *The Poincaré case.* Here f is assumed C^1 -close to Identity (without lying on a 1-parameter group).

Let \tilde{f} be a lift of f as in definition 1.11. We look at the graph $G(\tilde{f})$ in $\tilde{A} \times \tilde{A}$ endowed with the symplectic form $\omega \ominus \omega := dp' \wedge dq' - dp \wedge dq$. The area preserving property of f is equivalent to $G(\tilde{f})$ being a Lagrangian submanifold (i.e. $\omega \ominus \omega$ induces the null 2-form on $G(\tilde{f})$). Now we can make precise the closeness of f to Id :

$$(*) \quad \det(D\tilde{f} + Id) \neq 0.$$

This condition amounts to the transversality of $G(\tilde{f})$ to the anti-diagonal foliation.

Lemma 1.13. *Let $\Theta : \tilde{A} \times \tilde{A} \rightarrow T^*\tilde{A} \cong \mathbb{R}^2 \times \tilde{A}$ be defined by*

$$(p, q, p', q') \mapsto P_1 = q - q', \quad P_2 = p' - p, \quad Q_1 = \frac{p + p'}{2}, \quad Q_2 = \frac{q + q'}{2}.$$

Then:

1) Θ is a symplectic embedding,

2) the image of the diagonal Δ is $\{0\} \times \tilde{A}$,

3) Θ is equivariant with respect to the actions of \mathbb{Z} given by $(p, q, p', q') \mapsto (p+1, q, p'+1, q')$ and $(P_1, P_2, Q_1, Q_2) \mapsto (P_1, P_2, Q_1+1, Q_2)$.

Proof. By calculation we get $dP_1 \wedge dQ_1 + dP_2 \wedge dQ_2 = dp' \wedge dq' - dp \wedge dq$. The other points are obvious. \square

On $\Theta(G(\tilde{f}))$ we have:

$$\begin{aligned} (1) \quad & dP_1 \wedge dQ_1 + dP_2 \wedge dQ_2 = 0 \\ (2) \quad & P_1 = P_1(Q_1, Q_2), \quad P_2 = P_2(Q_1, Q_2). \end{aligned}$$

Condition (2) means that $\Theta(G(\tilde{f}))$ is the graph of some 1-form

$$\alpha(Q_1, Q_2) = P_1(Q_1, Q_2) dQ_1 + P_2(Q_1, Q_2) dQ_2,$$

and condition (1) means that $d\alpha = 0$. By the so-called Poincaré Lemma $\alpha = d\tilde{S}$ for some function $\tilde{S} : \tilde{A} \rightarrow \mathbb{R}$ and:

$$\begin{cases} q - q' = \partial_1 \tilde{S} \\ p' - p = \partial_2 \tilde{S}. \end{cases}$$

These last formulae make clear the following equivalence:

$$(p, q) \in \text{Fix } \tilde{f} \iff (p, q) \in \text{crit } \tilde{S}.$$

If $q = q' = -1$, then $Q_2 = \frac{q+q'}{2} = -1$, $P_1 = 0$ and $P_2 < 0$ according to the twist condition.

If $q = q' = +1$, then $Q_2 = +1$, $P_1 = 0$ and $P_2 > 0$. Therefore \tilde{S} is constant along each boundary component of \tilde{A} and the gradient points out.

Lemma 1.14. $\tilde{S} : \tilde{A} \rightarrow \mathbb{R}$ is 1-periodic in Q_1 .

Proof. The graph $G(\tilde{f})$ is \mathbb{Z} -invariant. Then according to 3) from lemma 1.13 the differential of \tilde{S} is 1-periodic in Q_1 . So it is sufficient to check the periodicity of \tilde{S} along $\{Q_2 = 1\}$. But there \tilde{S} is constant, *a fortiori* 1-periodic. \square

As a consequence \tilde{S} descends to the annulus $A = S^1 \times [-1, +1]$ hence defines a function $S : A \rightarrow \mathbb{R}$. Knowing the position of its gradient along the boundary and using the minimax argument (lemma 1.12), we see that S has at least two critical points. It finishes the proof that f has at least two fixed points.

Remark 1.15. A function S as above is called a *generating function* of the symplectic diffeomorphism f . The same idea does not work directly when f is not close enough to the Identity. But in [30] L. Hörmander introduced a generalization, called *generating phase* or *generating family* in Weinstein's terminology. With this tool M. Chaperon succeeded in giving a simple proof of the Arnold conjecture for the $2n$ -torus, previously proven by C. Conley and E. Zehnder [14]. As V.I. Arnold noticed in Appendix 9 from [1], the case $n = 1$ includes the C^1 -Poincaré-Birkhoff theorem.

2 Symplectic action on loops

2.1 Action functional in case of $(\mathbb{R}^{2n}, \omega_0 = dp \wedge dq)$

Let $H \in C^\infty(S^1 \times \mathbb{R}^{2n}, \mathbb{R})$ be a 1-periodic time dependent Hamiltonian and let H_t be its restriction to $\{t\} \times \mathbb{R}^{2n}$. Denote $\Lambda := C^\infty(S^1, \mathbb{R}^{2n})$. For $x \in \Lambda$, $x(t) = (p(t), q(t))$, we define the *action functional* by the following formula:

$$\Phi_H(x) = - \int_{S^1} x^* [p dq - H dt] = - \int_{S^1} [p(t)\dot{q}(t) - H(t, x(t))] dt.$$

If δx is a tangent vector to Λ at x (ie a vector field in \mathbb{R}^{2n} along x), the differential of Φ_H is well defined in the direction of δx and we have:

$$d\Phi_H(x) \cdot \delta x = - \int_{S^1} [\dot{q} \delta p + p \delta \dot{q} - dH_t(x) \cdot \delta x] dt,$$

After integration by parts we get:

$$\begin{aligned} d\Phi_H(x) \cdot \delta x &= - \int_{S^1} [\dot{q} \delta p - \dot{p} \delta q - dH_t(x) \cdot \delta x] dt \\ &= - \int_{S^1} [\omega_0(\delta x(t), \dot{x}(t)) - dH_t(x) \cdot \delta x] dt \\ &= + \int_{S^1} [\omega_0(\dot{x}, \delta x) + dH_t(x) \cdot \delta x] dt. \end{aligned}$$

This computation is basic in *the calculus of variations* and yields the following proposition:

Proposition 2.1. *The loop x is a critical point of $\Phi_H \Leftrightarrow \omega_0(\dot{x}, \cdot) + dH_t(x) = 0 \Leftrightarrow \dot{x} = X_{H_t}(x) \Leftrightarrow x$ is a 1-periodic solution of the Hamiltonian system from definition 1.4.*

2.2 Case of an autonomous C^2 -small Hamiltonian

Lemma 2.2. *Let $H : \mathbb{R}^{2n} \rightarrow \mathbb{R}$ be an autonomous Hamiltonian function satisfying $\|d\nabla H(x)\| < 2\pi$. Then any 1-periodic solution of the Hamiltonian system is constant and coincides with a critical point of H .*

Proof. The n -th Fourier coefficient of x is $c_n(x) = \int_0^1 e^{-2i\pi nt} x(t) dt$ and $c_n(\dot{x}) = -2i\pi n c_n(x)$. Since $c_0(\dot{x}) = 0$ from Parseval's formula we have:

$$\|\dot{x}\|_{L^2} \leq \frac{1}{2\pi} \|\ddot{x}\|_{L^2}.$$

But $\ddot{x} = J d\nabla H(x) \cdot \dot{x}$ implies $\|\ddot{x}\|_\infty < 2\pi \|\dot{x}\|_\infty$, which leads to a contradiction if x is not constant. \square

Corollary 2.3. *Let M be a closed symplectic manifold and $H : M \rightarrow \mathbb{R}$ an autonomous Hamiltonian. If H is C^2 -small enough, any 1-periodic solution of $\dot{x} = X_H(x)$ is constant.*

Proof. Let us equip M with a finite atlas of compact symplectic charts. This means a finite covering of M with the interiors of compact domains K_i on the neighbourhood of which symplectic charts are defined. If H is C^1 -small, every solution of $\dot{x} = X_H(x)$ has a small diameter and then it remains in one of the K_i 's. Here the C^2 norm is equivalent to the Euclidian norm of \dot{x} in the symplectic coordinates of K_i and the above lemma applies. \square

There are two consequences of the previous discussion.

1) If our goal is to estimate from below $\#\text{Per}(X_H)$, the number of 1-periodic solutions, independently of H , we see that we have to take only the contractible orbits into account.

2) In the case of Lemma 2.2 we have $\text{Per}(X_H) = \text{crit } H$ showing how the Morse theory comes into the problem of periodic orbits.

2.3 Action functional in the case of the $2n$ -torus

Let $\pi : \mathbb{R}^{2n} \rightarrow T^{2n}$ be the universal cover. Let us denote $\Lambda = C_c^\infty(S^1, T^{2n})$, the space of smooth contractible loops. For $x \in \Lambda$ we have $u : D^2 \rightarrow T^{2n}$, which is unique up to homotopy relatively to S^1 (because the 2nd homotopy group of T^{2n} is trivial), and $\tilde{x} : S^1 \rightarrow \mathbb{R}^{2n}$ such that x is the restriction of u to S^1 and $x = \pi \circ \tilde{x}$. Given a Hamiltonian H 1-periodic in time, the action integral is now defined by:

$$\Phi_H(x) = - \int_{S^1} \tilde{x}^*(p dq) + \int_{S^1} H(t, x) dt.$$

By the Stokes formula we also have:

$$\Phi_H(x) = - \int_{D^2} u^* \omega_0 + \int_{S^1} H(t, x) dt$$

where ω_0 denotes the standard symplectic structure on T^{2n} .

2.4 Homotopic assumption

In order to define the action integral for a more general symplectic manifold (M, ω) we will for simplicity make the following very strong assumption:

$$\langle \omega, \pi_2(M) \rangle = 0$$

ie: for any $V : S^2 \rightarrow M$, $\int_{S^2} v^* \omega = 0$. Of course this assumption is much stronger than *monotonicity*, the assumption made in Floer's work. If Λ still

denotes the space of smooth contractible loops of M , for $x \in \Lambda$ the following integral is well defined:

$$\Phi_H(x) = - \int_{D^2} u^* \omega + \int_{S^1} H(t, x) dt,$$

where $u : D^2 \rightarrow M$ extends $x : S^1 \rightarrow M$. The usual calculus of variations yields:

$$x \in \text{crit } \Phi_H \iff x \text{ is a 1-periodic solution of } \dot{x} = X_{H_t}.$$

In the sequel we will try to make a Morse theory for the action functional even though it is defined on an infinite dimensional space.

As a side-remark we observe that, when $\langle \omega, \pi_2(M) \rangle \neq 0$, the action functional becomes a multivalued function. Indeed for a given $x : S^1 \rightarrow M$ there are different $u : D^2 \rightarrow M$ providing different integrals $\int_{D^2} u^* \omega$. Only the differential of Φ_H is well defined. Hence the problem of counting the critical points deals with Morse-Novikov theory which is an analogue of Morse theory for closed 1-forms instead of real functions ([38], [28]). Actually there are other difficulties that appear simultaneously when the above homotopic assumption is not made, like the loss of compactness as it is stated in section 4.

2.5 L^2 -gradient of the action in the case of the standard torus

First let us recall J , the complex structure on \mathbb{C}^n , and its relation with the symplectic structure ω_0 and the Euclidean product. For $\xi, \eta \in \mathbb{C}^n$,

$$\omega_0(\xi, \eta) = - \langle \xi, J\eta \rangle.$$

Here \langle, \rangle stands for the Euclidean scalar product on \mathbb{R}^{2n} . Now if x is a smooth loop in T^{2n} and δx a variation of x we have:

$$\begin{aligned} d\Phi_H(x) \cdot \delta x &= \int \omega_0(\dot{x}, \delta x) - \omega_0(X_H(t, x), \delta x) dt \\ &= \int (- \langle \dot{x}, J\delta x \rangle + \langle X_H(t, x), J\delta x \rangle) dt \\ &= \int_{S^1} \langle J\dot{x} + \nabla H, \delta x \rangle dt \\ &= \langle J\dot{x} + \nabla H(t, x), \delta x \rangle_{L^2}. \end{aligned}$$

Therefore the L^2 -gradient of Φ_H is defined at $x \in \Lambda$ by:

$$\text{grad } \Phi_H(x) = J\dot{x} + \nabla H(t, x),$$

which is a vector field tangent to T^{2n} along the loop x .

Remark 2.4. One can check that the Cauchy problem associated with this gradient on Λ is ill-posed. In order to correct this fact the very seminal idea of Floer is to look at the gradient lines globally. As in Morse theory one looks at descending gradient lines. Such a line in Λ is a family of loops

$$\begin{aligned} u : \mathbb{R} \times S^1 &\rightarrow T^{2n} \\ (s, t) &\mapsto u(s, t) \end{aligned}$$

which is a solution of

$$\frac{\partial u}{\partial s} = -J \frac{\partial u}{\partial t} - \nabla H(t, u(s, t)).$$

This equation is nothing but a Cauchy-Riemann equation perturbed by a non-linear term of order 0:

$$\bar{\partial}u + \nabla H(t, u) = 0$$

where $\bar{\partial}u = \frac{\partial u}{\partial s} + J \frac{\partial u}{\partial t}$.

Proposition 2.5. *If $u : \mathbb{R} \times S^1 \rightarrow T^{2n}$ is a solution of the perturbed Cauchy-Riemann equation, then we have:*

$$\begin{aligned} \Phi_H(u_{s_1}) - \Phi_H(u_{s_0}) &= - \int_{s_0}^{s_1} \|\text{grad } \Phi_H\|_{L^2}^2 ds dt \\ &= - \int_{s_0}^{s_1} \int_{S^1} \|J \frac{\partial u}{\partial t} + \nabla H(t, u)\|^2 ds dt. \end{aligned}$$

2.6 L^2 -gradient of the action for a general symplectic manifold

This requires the choice of a compatible almost complex structure.

Definition 2.6. *An almost complex structure J on the manifold M is a linear isomorphism of the tangent space TM such that $J^2 = -I$. If ω is a symplectic structure, J is said ω -compatible if $\omega(\xi, J\eta)$ is a scalar product.*

Proposition 2.7. *The space $\mathcal{J}(\omega)$ of the almost complex structures ω -compatible is contractible (in particular non-empty).*

Proof. This is based on the *polar decomposition* of a linear isomorphism of $T_a M$. Let g be an auxiliary Riemannian metric on M . We have $\omega(\xi, \eta) = g(\xi, \Omega\eta)$ and ${}^t\Omega = -\Omega$ is non-singular. There exists a unique symmetric isomorphism S , positive definite, such that $S^2 = \Omega {}^t\Omega$. Let us put $J = \Omega^{-1}S$. One easily checks:

- J is g -unitary,
- $J^{-1}SJ$ is symmetric positive definite and it is a square root of $\Omega {}^t\Omega$, so $J^{-1}SJ = S$,

- $S^2 J^{-2} = (S J^{-1})^2 = \Omega^2 = -S^2$, and then $J^2 = -I$,
- $\omega(\xi, J\eta) = g(\xi, \Omega J\eta) = g(\xi, S\eta)$ is a scalar product.

Finally the 1-1 correspondence $g \leftrightarrow J$ answers the question. \square

For the sequel we choose such a J and we define the Riemannian metric

$$\langle \xi, \eta \rangle = \omega(\xi, J\eta).$$

If ∇H is the gradient of H (t being kept fixed) with respect to this metric, the Hamiltonian system reads

$$\dot{x} = J(x) \nabla H(t, x).$$

By a similar calculation to the case of $M = T^{2n}$, the L^2 -gradient of the action functional is given by

$$\text{grad } \Phi_H(x) = J(x)\dot{x} + \nabla H(t, x)$$

and the equation of the gradient lines of Φ_H reads:

$$\frac{\partial u}{\partial s} + J(u) \frac{\partial u}{\partial t} + \nabla H(t, u) = 0.$$

3 A brief survey of Morse theory

3.1 Dynamical systems point of view

Definition 3.1. A function $H : M \rightarrow \mathbb{R}$ is called a *Morse function* if for any $a \in \text{crit } H$, $d^2 H(a) : T_a M \times T_a M \rightarrow \mathbb{R}$ is non-degenerate. The Morse index $\text{ind}_a H = \text{ind}(a)$ is the maximal dimension of a subspace in $T_a M$ on which $d^2 H(a)$ is negative definite.

As a consequence of the implicit function theorem the critical points of a Morse function are isolated.

Theorem 3.2. (Thom [54]) *When M is closed the set of Morse functions is open and dense in $C^\infty(M, \mathbb{R})$. In particular to be a Morse function is a generic property in the sense of Baire.*

Lemma 3.3. (Morse's Lemma) *If a is a Morse critical point of H there exist local coordinates around a such that*

$$H(x) = H(a) - x_1^2 - \dots - x_i^2 + x_{i+1}^2 + \dots$$

where $i = \text{ind}_a H$.

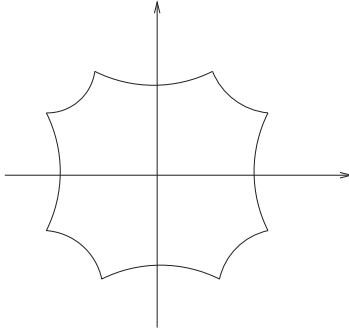


Figure 3

Such a chart is called a *Morse chart*. It contains Morse models an example of which is drawn in figure 3. A *Morse model* is a compact domain in a Morse chart whose frontier consists of three parts: the lower (resp. upper) side which is contained in a regular level set of H below (resp. above) the critical value, and the lateral side which is a union of gradient segments joining both other sides. We will choose a Riemannian metric on M which coincides with the Euclidean metric in some Morse charts around all the critical points. We look at the descending gradient lines, solutions of

$$\dot{x} = -\nabla H(x).$$

Let χ_t be its flow. The unstable (resp. stable) manifold of a is

$$W^u(a) \text{ (resp. } W^s(a)) := \{x \in M \mid \chi_t(x) \xrightarrow[t \rightarrow \mp\infty]{} a\}.$$

According to Morse's Lemma and the choice of the Riemannian metric, $W^u(a) \cong \mathbb{R}^{\text{ind}_a H}$.

Theorem 3.4. (Smale [52]) *A generic metric (among those which coincide with a given metric near the critical points) meets the Smale transversality condition: the stable and unstable manifolds are mutually transversal.*

Comment. The reason why the extra (closed) condition imposed to the metrics does not imply loss of genericity is the following. Starting from a metric g_0 , in order to make the stable and unstable manifolds mutually transversal it is sufficient to modify g_0 near a family of regular level sets made of one between each pair of consecutive critical level sets.

The next theorem does not require the Smale transversality condition. On a closed manifold it is elementary to see that, for any initial condition, the corresponding solution of the gradient equation tends to a critical point when $t \rightarrow \pm\infty$. Let \mathcal{M} denote the set of solutions (or *gradient lines*) and, for $a, b \in \text{crit } H$, let $\mathcal{M}(a, b)$ denote the set of gradient lines descending from a to b :

$$\mathcal{M} = \bigcup_{a, b \in \text{crit } H} \mathcal{M}(a, b),$$

including $\mathcal{M}(a, a)$ which is nothing but the constant a . Here a gradient line means a *parametrized gradient line*. We have a natural evaluation map

$$\begin{aligned} ev : \mathcal{M} &\rightarrow M \\ \gamma &\mapsto \gamma(0) \end{aligned}$$

and a right action of \mathbb{R} on \mathcal{M} by translation of time:

$$(\gamma \in \mathcal{M}, s \in \mathbb{R}) \mapsto (\gamma \cdot s)(t) := \gamma(s + t).$$

An equivalence class $\bar{\gamma}$ under this action is an unparametrized gradient line or a *connecting orbit*.

Theorem 3.5. *If M is a closed manifold and \mathcal{M} is endowed with the $C_{loc}^\infty(\mathbb{R}, M)$ topology, then we have the following:*

- 1) \mathcal{M} is a compact set and $ev : \mathcal{M} \rightarrow M$ is a homeomorphism.
- 2) If $\gamma_n \rightarrow \gamma$ is a converging sequence in \mathcal{M} , with $\gamma_n \in \mathcal{M}(a, b)$ for all n , then there exist a subsequence (still denoted by γ_n) and :

- critical points $a_0 = a, a_1, \dots, a_\ell, a_{\ell+1} = b$,
- for each $k \in \{0, \dots, \ell\}$, a sequence of real numbers $(s_n^k)_n$ and a gradient line $\gamma^k \in \mathcal{M}(a_k, a_{k+1})$

such that $\gamma_n \cdot s_n^k \rightarrow \gamma^k$ for every k .

- 3) The limit γ is one of the γ^k 's (up to a finite shift of the variable) if γ is not constant, or one of the a_k 's if it is.

The $(\ell + 1)$ -uple of connecting orbits $(\bar{\gamma}^0, \dots, \bar{\gamma}^\ell)$ is called a *broken connecting orbit* from a to b .

Proof. 1) This is just the Cauchy-Lipschitz theorem.

2) There is a simple proof of this point using Morse models. But we take care to give a proof which also works in infinite dimension. We look at a regular value h_1 of H less than $H(a)$ so that $[h_1, H(a)[$ contains no critical values. Let s_n^1 be real numbers such that $(H \circ \gamma_n)(s_n^1) = h_1$. Taking a subsequence if necessary, $\gamma_n \cdot s_n^1 \rightarrow \gamma^1$ in \mathcal{M} and γ^1 belongs to some $\mathcal{M}(a_0, a_1)$. We leave as an exercise to prove that $a_0 = a$ (argue as below). We have:

$$H(a_1) = \inf_T \lim_n H(\gamma_n(s_n^1 + T)).$$

For the next step we take a regular value h_2 of H with no critical values in $[h_2, H(a_1)[$. Let s_n^2 be real numbers such that $(H \circ \gamma_n)(s_n^2) = h_2$. Up to taking a suitable subsequence, we have:

$$\gamma_n \cdot s_n^2 \rightarrow \gamma^2 \in \mathcal{M}(a'_1, a_2).$$

Certainly $s_n^2 - s_n^1$ is unbounded, if not $\gamma^2 = \gamma^1$ up to a finite shift of the variable, and

$$h_2 = H(\gamma_n(s_n^2)) = H(\gamma^2(0)) > H(a_1),$$

contradicting $h_2 < H(a_1)$. Hence, for every $T > 0$ and n big enough we have

$$H(\gamma_n(s_n^2 - T)) < H(\gamma_n(s_n^1 + T))$$

which yields $H(a'_1) \leq H(a_1)$. Since $[h_2, H(a_1)[$ contains no critical values, we have $H(a'_1) = H(a_1)$.

Let us prove that $a'_1 = a_1$. If not, there exist small disjoint open balls B and B' centered at a_1 and a'_1 whose closure contains no other critical points of H other than their centers. There exists τ such that $\gamma^1(t) \in B$ and $\gamma^2(-t) \in B'$ for every $t \geq \tau$. For n big enough, $\gamma_n(s_n^1 + \tau) \in B$ and $\gamma_n(s_n^2 - \tau) \in B'$. So, for every n , there exists $\theta_n \in]s_n^1 + \tau, s_n^2 - \tau[$ such that $\gamma_n(\theta_n) \in S$, where S is the frontier of B . Taking a suitable subsequence we have $\gamma_n \cdot \theta_n \rightarrow \gamma' \in \mathcal{M}$. The sequence $\theta_n - s_n^1$ is unbounded. If not, as above, $\gamma' = \gamma^1$ up to a finite shift; but $\gamma^1(] \tau, +\infty[)$ does not meet S ; contradiction. In the same way $s_n^2 - \theta_n$ is unbounded and the limit values of $H \circ \gamma'$ satisfy

$$H(a_1) \geq \sup H \circ \gamma'(\mathbb{R}) \geq \inf H \circ \gamma'(\mathbb{R}) \geq H(a'_1) = H(a_1).$$

Therefore γ' is constant and $\gamma'(0) \in S$ is a critical point of H . Contradiction; so $a'_1 = a_1$. The other steps go in the same way with a slight modification for $a_{\ell+1} = b$.

3) Consider the case where $H(\gamma(0)) \in]H(a_{k+1}), H(a_k)[$ for some k . Suppose s_n^k is unbounded; let us say $s_n^k \rightarrow +\infty$ (the other cases are left to the reader). Hence

$$\lim H(\gamma_n(0)) \geq \sup H(\gamma^k(\mathbb{R})) = H(a_k).$$

So s_n^k is bounded and $\gamma = \gamma^k$ up to a finite shift of the variable.

If now $H(\gamma(0)) = H(a_k)$, then necessarily $s_n^k \rightarrow +\infty$ and $s_n^{k-1} \rightarrow -\infty$ and we conclude that $H \circ \gamma$ is constant. But $\bar{\gamma}$ is a gradient line. So it is a critical point c . As in the end of 2) $c \neq a_k$ leads to a contradiction with the fact that the critical points are isolated. \square

From now on, we assume the Smale transversality condition as in theorem 3.4. Unfortunately in spite of the above theorem the picture of the closure of $\mathcal{M}(a, b)$ is still a little difficult to describe, (but it is done in [33]). If we look at the quotient space \mathcal{M}/\mathbb{R} of the connecting orbits, the picture is horrible because this space is not Hausdorff. But there is an *abstract* compactification of $\mathcal{M}(a, b)/\mathbb{R}$

which has a very nice structure as it is explained in the statement below, which is a direct consequence of the transversality condition and of theorem 3.5.

Corollary 3.6.

- 1) $\mathcal{M}(a, b)/\mathbb{R}$ is a manifold of dimension $\text{ind}(a) - \text{ind}(b) - 1$. In particular if $\text{ind}(b) \geq \text{ind}(a)$, $\mathcal{M}(a, b)/\mathbb{R}$ is empty.
- 2) There is a compact manifold with boundary and corners $\widehat{\mathcal{M}(a, b)}/\mathbb{R}$ formed by adding all broken orbits from a to b .

Part 1) is obvious since $W^u(a) \cap W^s(b)$ is a manifold whose dimension is $\text{ind}(a) - \text{ind}(b)$, and the action of \mathbb{R} on it is free and proper. A proof of part 2) can be read in [31] or [46]. We will detail below the easy particular case when the difference of indices is 1 or 2.

Corollary 3.7.

- 1) If $\text{ind}(a) - \text{ind}(b) = 1$, $\mathcal{M}(a, b)/\mathbb{R}$ is made of finitely many points.
- 2) If $\text{ind}(a) - \text{ind}(b) = 2$, $\mathcal{M}(a, b)/\mathbb{R}$ is made of a finite number of closed intervals or circles.

Proof. 1) Due to the indices there are no broken connecting orbits from a to b and we have to show $\mathcal{M}(a, b)/\mathbb{R}$ is already a compact set. Let $(\overline{\gamma}_n)_n$ be a sequence in the quotient space. It can be lifted to $\mathcal{M}(a, b)$ as a sequence $(\gamma_n)_n$ such that $\gamma_n(0)$ belongs to a given regular level set. Then, according to theorem 3.5, every cluster point of this sequence in \mathcal{M} belongs to $\mathcal{M}(a, b)$. So $\mathcal{M}(a, b)/\mathbb{R}$ is compact. Being a 0-dimensional manifold it is a finite set.

2) Here a broken connecting orbit is made of exactly two pieces $(\overline{\gamma}_0, \overline{\gamma}_1)$. Let c be the intermediate critical point. Let k be its Morse index and n the ambient dimension. We look at a Morse model about c . The upper side V_+ , contained in $\{H = H(c) + \varepsilon\}$, is diffeomorphic to $D^k \times S^{n-k-1}$ and $W^u(a)$ traces a k -manifold Σ_+ on V_+ which meets $0 \times S^{n-k-1}$, the trace of $W^s(c)$, transversally in a single point p corresponding to $\overline{\gamma}_0$. The lower side V_- of the Morse model, contained in $\{H = H(c) - \varepsilon\}$, is diffeomorphic to $S^{k-1} \times D^{n-k}$. The gradient lines yield a diffeomorphism

$$G : V_+ \setminus (0 \times S^{n-k-1}) \rightarrow V_- \setminus (S^{k-1} \times 0).$$

If we use *polar coordinates*, $\varphi \in S^{k-1}$, $\psi \in S^{n-k-1}$, $r \in]0, 1]$, a point of V_+ is given by $((\varphi, r), \psi)$, a point of V_- is given by $(\varphi, (\psi, r))$ and, when $r \neq 0$,

$$G(\varphi, r, \psi) = (\varphi, \psi, r).$$

With the help of this model it is easy to see that $G(\Sigma_+ \setminus p)$ is a punctured manifold which compactifies as a manifold with boundary Σ_- when adding $S^{k-1} \times 0$ as a boundary. Now the trace of $W^s(b)$ on V_- is a $(n-k)$ -dimensional manifold which intersects $S^{k-1} \times 0$ transversally in exactly one point corresponding to $\overline{\gamma}_1$.

Therefore its intersection with Σ_- is a 1-dimensional manifold with boundary. As the interior points of this arc represent gradient lines descending directly from a to b , this arc provides $\widehat{\mathcal{M}(a, b)}/\mathbb{R}$ with a chart of a 1-manifold with boundary near $(\bar{\gamma}_0, \bar{\gamma}_1)$. \square

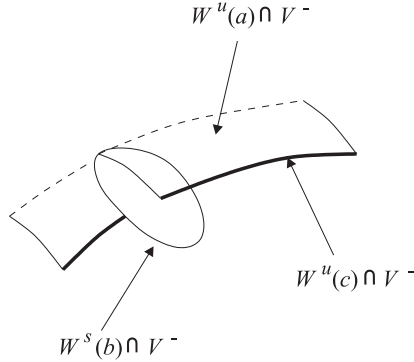


Figure 4: Traces of stable/unstable manifolds in the lower side V_- of the Morse model.

3.2 The Morse-Smale-Witten complex

This complex is also called the Thom-Smale complex because of the note by R. Thom [53] which is the very beginning of this story. But the idea of transversality was missing and the right statement was given by S. Smale in [52]. More recently E. Witten rediscovered this complex by another approach.

For simplicity we do not take the orientations into account. Therefore all the coefficients are in $\mathbb{Z}/2$. Here H will be a Morse function on the closed manifold M and g will be a generic Riemannian metric as in theorem 3.4.

Notation 3.8. $C_k(H, g)$ is the $\mathbb{Z}/2$ -vector space generated by $\text{crit}_k(H)$, the set of critical points of index k . For $k \notin \{0, \dots, \dim M\}$, $C_k = 0$.

If $a \in \text{crit}_k(H)$ and $b \in \text{crit}_{k-1}(H)$, $n(a, b)$ denotes the number of connecting orbits from a to b .

The boundary operator $\partial_k : C_k(H, g) \rightarrow C_{k-1}(H, g)$ is defined by

$$\partial_k(a) = \sum_{b \in \text{crit}_{k-1}H} n(a, b) b.$$

Theorem 3.9.

- 1) $\partial_k \circ \partial_{k+1} = 0$.
- 2) $H_k(H, g) := \frac{\ker \partial_k}{\text{im } \partial_{k+1}}$ is a $\mathbb{Z}/2$ -vector space of finite dimension.
- 3) (Poincaré duality) $H_k(H, g)$ and $H_{n-k}(-H, g)$ are dual one to the other in a canonical way.

Proof.

1) If $b \in \text{crit}_{k-1}(H)$ contributes to $\partial_k \circ \partial_{k+1}(a)$, this means there exists a broken orbit (a, c, b) where c is some critical point of index k . It corresponds to an endpoint of some interval in $\mathcal{M}(\widehat{a, b})/\mathbb{R}$. Then the other endpoint of this interval cancels the contribution mod 2 of c to the coefficient of b in $\partial_k \circ \partial_{k+1}(a)$.

2) This is clear because $C_k(H, g)$ is finitely generated.

3) If a is a critical point of H of index k , it is a critical point of $-H$ of index $n - k$. If γ is a connecting orbit descending from a to b for H , then γ^* , that is γ traversed in the opposite way, is a connecting orbit descending from b to a for $-H$. Let us consider the pairing

$$C_k(H, g) \otimes C_{n-k}(-H, g) \rightarrow \mathbb{Z}/2$$

defined at the level of the critical points by $\langle a, b \rangle = 1$ if and only if $a = b$. According to the definition of the boundary operator, the one from $C_*(H, g)$ is the transpose of the one from $C_*(-H, g)$. \square

Theorem 3.10. (Continuation theorem)

- 1) If (H_λ, g_λ) , $\lambda \in [0, 1]$, is a generic path of pairs, each being formed with a function and a Riemannian metric, it induces a morphism

$$\Phi_* : H_*(H_0, g_0) \rightarrow H_*(H_1, g_1).$$

- 2) Φ_* does not depend on the chosen path.

- 3) (Naturality) If we have a diagram

$$H_*(H_0, g_0) \xrightarrow{\Phi_*^{10}} H_*(H_1, g_1) \xrightarrow{\Phi_*^{21}} H_*(H_2, g_2) \xrightarrow{\Phi_*^{02}} H_*(H_0, g_0),$$

then $\Phi_*^{02} \circ \Phi_*^{21} \circ \Phi_*^{10} = \text{Id}$. In particular Φ_*^{10} is an isomorphism whose inverse is Φ_*^{01} .

This statement allows us to define $H_*(M, \mathbb{Z}/2)$ as the “common value” of all homologies $H_*(H, g)$. Notice also that the functions from the path in 1) are allowed to be non-Morse for a finite number of values of the parameter.

Proof. We only pay attention to the first part, whose proof yields easily the other parts. We will show that the given path allows us to define a morphism of complexes

$$\Phi : C_*(H_0, g_0) \rightarrow C_*(H_1, g_1),$$

which means a morphism commuting with the boundary operators. It is obtained in the following way according to an idea of Floer. One integrates the differential equation

$$\frac{dx}{d\lambda} = -\nabla H_\lambda.$$

For a generic path, there are finitely many connecting orbits from critical points of H_0 to critical points of H_1 whose indices are the same.

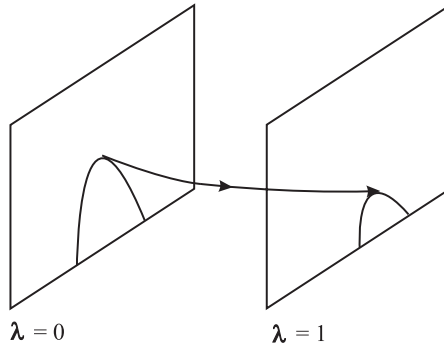


Figure 5: A gradient line of K descending from $M \times \{0\}$

- A way to see this is to introduce a smooth decreasing function $\rho(\lambda)$ whose derivative vanishes exactly at the endpoints of $[0, 1]$ to order 1. Let $K : M \times [0, 1] \rightarrow \mathbb{R}$ be defined by $K(x, \lambda) = H_\lambda(x) + \rho(\lambda)$. If ρ decreases strongly enough, dK vanishes only at the critical points of H_i in $M \times \{i\}$ for $i = 0, 1$. The gradient lines on the product yields the desired connecting orbits. Observe that a critical point a of H_0 satisfies $\text{ind}_a K = \text{ind}_a H_0 + 1$. But a critical point a of H_1 satisfies $\text{ind}_a K = \text{ind}_a H_1$.

Now, for $a \in \text{crit}_k(H_0)$, we set

$$\Phi(a) = \sum_b n(a, b) b$$

where the sum is taken over $\text{crit}_k(H_1)$, $n(a, b)$ being the number of connecting orbits from a to b mod 2.

The fact that Φ commutes with ∂ can be deduced from $\partial \circ \partial = 0$ when we look at the above mentioned function K on the product $M \times [0, 1]$. \square

4 Compactness

4.1 Finite energy

We recall the situation: (M, ω) is a closed symplectic manifold, $H : S^1 \times M \rightarrow \mathbb{R}$ is a Hamiltonian 1-periodic in time, J is an almost complex structure ω -compatible and $\langle \xi, \eta \rangle = \omega(\xi, J\eta)$ is the associated Riemannian metric. We have an action functional Φ_H defined on the space Λ of contractible smooth loops in M :

$$\Phi_H = - \int_{D^2} v^* \omega + \int_{S^1} H(t, x) dt,$$

where $x : S^1 \rightarrow M$, $v : D^2 \rightarrow M$ and $v|_{\partial D^2} = x$. It is well defined under the homotopic assumption

$$\langle \omega, \pi_2(M) \rangle = 0.$$

The critical points of Φ_H are the 1-periodic solutions of the Hamilton differential equation

$$\dot{x} = J(x) \nabla H(t, x).$$

The descending L^2 -gradient lines of Φ_H are the solutions of the perturbed Cauchy-Riemann equation

$$(*) \quad \frac{\partial u}{\partial s} + J(u) \frac{\partial u}{\partial t} + \nabla H(t, u) = 0.$$

We define the *energy* of a gradient line u by

$$\begin{aligned} E(u) &:= - \int_{-\infty}^{+\infty} \frac{d}{ds} \Phi_H(u(s)) ds \\ &= \frac{1}{2} \int_{\mathbb{R} \times S^1} \left(\left\| \frac{\partial u}{\partial s} \right\|^2 + \left\| J(u) \frac{\partial u}{\partial t} + \nabla H(t, u) \right\|^2 \right) ds dt. \end{aligned}$$

Here $u(s)$ means $u|_{\{s\} \times S^1}$. Of course both terms from the last integrand are equal. We introduce the set of gradient lines whose energy is finite:

$$\mathcal{M} = \{u : \mathbb{R} \times S^1 \rightarrow M \mid \text{smooth, contractible, solution of } (*), E(u) < \infty\}.$$

Proposition 4.1.

- 1) If $u \in \mathcal{M}$, $\Phi_H(u(s))$ tends to a critical value $\Phi(u_{\pm})$ when s goes to $\pm\infty$, where u_{\pm} is a 1-periodic solution of the Hamilton equation and there is at least one accumulating point of $\{u(s) \mid s \rightarrow \pm\infty\}$ in Λ .
- 2) There exists a constant $C > 0$ such that $\Phi_H(\mathcal{M}) \subset [-C, C]$.
- 3) For $u \in \mathcal{M}$, $E(u) \leq 2C$.

Proof.

1) For $u \in \mathcal{M}$ we have:

$$\int_{-\infty}^{+\infty} \int_{S^1} \|\dot{u} - J(u)\nabla H(t, u)\|^2 ds dt < +\infty$$

where $\dot{u} = \frac{\partial u}{\partial t}$. Then there exists a sequence $s_k \rightarrow +\infty$ (or $-\infty$) such that

$$(**) \quad \|\dot{u}_k - J(u_k)\nabla H(t, u_k)\|_{L^2} \rightarrow 0$$

where $u_k := u(s_k)$. In order to speak of L^2 we use an embedding of M into some \mathbb{R}^N .

As ∇H is bounded (compactness of M), there exists $B > 0$ such that, for every k , $\|\dot{u}_k\|_{L^2} \leq B$. This bound implies that the family of the u_k 's is equicontinuous. Indeed:

$$u_k(t_1) - u_k(t_0) = \int_{t_0}^{t_1} \dot{u}_k = \int_{S^1} \mathbf{1}_{[t_0, t_1]} \dot{u}_k \leq [t_1 - t_0]^{1/2} \|\dot{u}_k\|_{L^2} \leq B[t_1 - t_0]^{1/2}.$$

Thus the family has a compact closure in $C^0(S^1, M)$ (Arzelà-Ascoli Theorem). Taking a subsequence if necessary $u_k \rightarrow u_+$ (or u_-) in C^0 . Due to (**), \dot{u}_k is a Cauchy sequence in L^2 converging to the continuous function $v = J(u_+)\nabla H(t, u_+)$. By the dominated convergence theorem,

$$u_+(t) - u_+(0) = \int_0^t v(\theta) d\theta.$$

Therefore u_+ is C^1 and a solution of the Hamilton equation. Then it is C^∞ . Of course $\Phi_H(u_k) \rightarrow \Phi_H(u_+)$ with $k \rightarrow \infty$. But, as $\Phi_H(u(s))$ is monotone, we get

$$\Phi_H(u(s)) \xrightarrow{s \rightarrow +\infty} \Phi_H(u_+).$$

2) The Arzelà-Ascoli Theorem implies easily that the set of the 1-periodic solutions of the Hamilton equation is compact in $C^\infty(S^1, M)$. Then the action functional is bounded on this set.

3) This is clear because $E(u) = \Phi_H(u_-) - \Phi_H(u_+)$. □

Corollary 4.2. (Exercise, [29] p. 236) \mathcal{M} is closed in $C_{loc}^\infty(\mathbb{R} \times S^1, M)$.

Remark 4.3. As pointed out by David Hermann, the arguments used in the above proposition imply that the action functional meets the so-called Palais-Smale condition: if $(x_n)_n$ is a sequence in Λ such that $\text{grad } \Phi(x_n)$ tends to 0 in L^2 , then there is a subsequence converging to a critical point of $\text{grad } \Phi$. But according to remark 2.4, this observation is of no help for the compactness of \mathcal{M} .

4.2 Compactness theorem

Theorem 4.4. *Under the assumption $\langle \omega, \pi_2(M) \rangle = 0$ (which allows us to define the action functional) we have the following:*

- 1) *There exists $A > 0$ such that, for every $u \in \mathcal{M}$ and $(s, t) \in \mathbb{R} \times S^1$, $\|\nabla u(s, t)\| \leq A$.*
- 2) *\mathcal{M} is compact in $C_{loc}^\infty(\mathbb{R} \times S^1, M)$.*

Proof. 1 \implies 2

By 1) and the Arzelà-Ascoli theorem, \mathcal{M} has a compact closure in $C_{loc}^0(\mathbb{R} \times S^1, M)$. Moreover the Cauchy-Riemann equation is elliptic. Then the so-called elliptic regularity implies that the solutions are C^∞ , any C^0 -limit of solutions is a solution and both C^0 and C^∞ topologies coincide on \mathcal{M} . Hence \mathcal{M} is C^0 -compact and then C^∞ -compact.

Proof of 1). Here it is more convenient to think of $u \in \mathcal{M}$ as defined on \mathbb{R}^2 instead of $\mathbb{R} \times S^1$, in particular when a rescaling is performed. If 1) is not true, there exist sequences $u_k \in \mathcal{M}$ and $(s_k, t_k) \in \mathbb{R}^2$ such that $\|\nabla u_k(s_k, t_k)\| \rightarrow \infty$. Let $\varepsilon_k > 0$ such that $\varepsilon_k \rightarrow 0$ and $\varepsilon_k \|\nabla u_k(s_k, t_k)\| \rightarrow \infty$.

Lemma 4.5. (Ekeland [18]) *Let $g : X \rightarrow \mathbb{R}_+$ be a continuous function on a complete metric space. Given $x_0 \in X$ and $\varepsilon_0 > 0$, there exist $x^* \in X$ and $0 < \varepsilon^* \leq \varepsilon_0$ such that:*

- 1) $d(x^*, x_0) \leq 2\varepsilon_0$,
- 2) $\varepsilon^* g(x^*) \geq \varepsilon_0 g(x_0)$,
- 3) $g(x) \leq 2g(x^*)$ for every $x \in B(x^*, \varepsilon^*)$.

Proof. If (x_0, ε_0) satisfies $g(x) \leq 2g(x_0)$ for every $x \in B(x_0, \varepsilon_0)$ one takes $x^* = x_0$ and $\varepsilon^* = \varepsilon_0$. If not we choose $x_1 \in B(x_0, \varepsilon_0)$ with $g(x_1) > 2g(x_0)$ and $\varepsilon_1 = \varepsilon_0/2$ and so on. The complete proof is left to the reader. \square

Let us apply this lemma to $g = \|\nabla u_k(s, t)\|$. Then there exist other ε_k and (s_k, t_k) such that:

- $\varepsilon_k \|\nabla u_k(s_k, t_k)\| \rightarrow \infty$ with k ,
- $2\|\nabla u_k(s_k, t_k)\| \geq \|\nabla u_k(s, t)\|$ for every $(s, t) \in B((s_k, t_k), \varepsilon_k)$.

The sequel of the proof consists in some rescaling:

$$v_k = u_k \left((s_k, t_k) + \frac{(s, t)}{R_k} \right)$$

where $R_k = \|\nabla u_k(s_k, t_k)\|$. The formula makes sense when u_k is thought of as defined on \mathbb{R}^2 , as we said above. We have some obvious properties:

- $\|\nabla v_k(0, 0)\| = 1$,

- $\|\nabla v_k(s, t)\| \leq 2$ on the “big” disk $D_k = B(\underline{0}, \varepsilon_k R_k)$ centered at the origin,
- $\frac{\partial v_k}{\partial s} + J(v_k) \frac{\partial v_k}{\partial t} + \frac{1}{R_k} \nabla H\left(t_k + \frac{t}{R_k}, v_k\right) = 0$.

Lemma 4.6. *For k big enough, we have*

$$\int_{D_k} \left\| \frac{\partial v_k}{\partial s} \right\|^2 + \left\| \frac{\partial v_k}{\partial t} \right\|^2 \leq 7C.$$

Proof. Let $B_k := B((s_k, t_k), \varepsilon_k)$. It is a small disk in \mathbb{R}^2 or in the cylinder as well. The first identity below is nothing but the change of variables in integrals and the rest is obvious.

$$\begin{aligned} \int_{D_k} \|\nabla v_k\|^2 &= \int_{B_k} \|\nabla u_k\|^2 \\ &\leq \int_{B_k} \left\| \frac{\partial u_k}{\partial s} \right\|^2 + \left\| \frac{\partial u_k}{\partial t} - X_H(t, u_k) + X_H(t, u_k) \right\|^2 \\ &\leq \int_{B_k} \left[\left\| \frac{\partial u_k}{\partial s} \right\|^2 + \left\| \frac{\partial u_k}{\partial t} - X_H \right\|^2 \right] + \|X_H\|^2 + 2 \left\| \frac{\partial u_k}{\partial t} - X_H \right\| \cdot \|X_H\| \\ &\leq \int_{B_k} \left[\left\| \frac{\partial u_k}{\partial s} \right\|^2 + 2 \left\| \frac{\partial u_k}{\partial t} - X_H \right\|^2 \right] + 2\|X_H\|^2 \quad (2ab \leq a^2 + b^2) \\ &\leq 3E(u_k) + 2 \int_{B_k} \|X_H\|^2 \leq 6C + 2 \int_{B_k} \|X_H\|^2. \end{aligned}$$

The last inequality follows from proposition 4.1 and the last integral converges to 0 with $k \rightarrow \infty$.

□

Lemma 4.7. (The bubbling lemma.) *Taking a subsequence if necessary, $v_k \rightarrow v$ in $C_{loc}^\infty(\mathbb{R}^2, M)$ with the following properties:*

- $\|\nabla v(0, 0)\| = 1$ (in particular v is not constant),
- $\|\nabla v(s, t)\| \leq 2$ for every $(s, t) \in \mathbb{R}^2$.
- $\frac{\partial v}{\partial s} + J(v) \frac{\partial v}{\partial t} = 0$ which means that v is J -holomorphic,

$$\bullet \int_{\mathbb{R}^2} \|\nabla v\|^2 < \infty.$$

Proof. Using the previous lemma, the bubbling lemma follows from usual Sobolev techniques as for $1 \Rightarrow 2$. \square

In order to finish the proof of theorem 4.4 let us show how the bubbling lemma leads to a contradiction. We have

$$\int_{\mathbb{R}^2} v^* \omega = \int_{\mathbb{R}^2} \omega \left(\frac{\partial v}{\partial s}, \frac{\partial v}{\partial t} \right) ds dt = \int \omega \left(\frac{\partial v}{\partial s}, J(v) \frac{\partial v}{\partial s} \right) = \int \left\| \frac{\partial v}{\partial s} \right\|^2 < \infty.$$

Since $\frac{\partial v}{\partial s}$ and $\frac{\partial v}{\partial t}$ are isometric and orthogonal, the last integral is the area of v . Let us put $A(r) := \int_{B(r)} v^* \omega$, the area of $v(B(r))$, and $\ell(r)$ the length of $v(\partial B(r))$. We know that $A(r)$ is a bounded primitive of $\ell(r)$. Then there exists a sequence $r_k \rightarrow \infty$ such that $\ell(r_k) \rightarrow 0$ (we see the bubble).

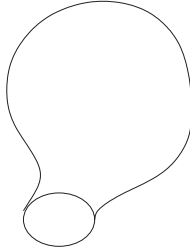


Figure 6: J -holomorphic disk with a small boundary.

For k large, $v(\partial B(r_k))$ has a small diameter and lies entirely in one of the charts from a finite atlas of M . Taking a subsequence if necessary all the $v(\partial B(r_k))$'s lie in the same chart U where ω is exact, $\omega = d\lambda$. We have

$$\left| \int_{v(\partial B(r_k))} \lambda \right| \leq \text{const. } \ell(r_k)$$

which converges to 0 with $k \rightarrow \infty$. There is a disk Δ_k in U (maybe with singularities) whose boundary is $v(\partial B(r_k))$ and Stokes formula applies and yields

$$\int_{\Delta_k} \omega \rightarrow 0.$$

We remember the homotopic assumption $\langle \omega, \pi_2(M) \rangle = 0$. So we have

$$\int_{v(B(r_k))} \omega = \left| \int_{\Delta_k} \omega \right|.$$

The right-hand side converges to 0. But the left-hand side converges to the area of v which is positive because v is a non constant J -holomorphic curve. This finishes the proof of the compactness theorem. \square

4.3 Non-degeneracy condition

In the sequel we shall assume the following generic condition: *Let φ_1 be the time 1 map of the Hamiltonian flow; for every fixed point a of φ_1 , 1 is not an eigenvalue of $D\varphi_1(a)$.* By the implicit function theorem a is an isolated fixed point of φ_1 . In that case we say that the 1-periodic solutions of the Hamilton equation (or of X_{H_t}) are *non-degenerate*.

Be careful! The non-degeneracy condition is never satisfied for a Hamiltonian independent of time if there are non-constant 1-periodic orbits.

Looking at $u \in \mathcal{M}$ as a curve in Λ , the space of smooth contractible loops in M , we are interested in its cluster points. For this purpose it is useful to introduce the *evaluation map*

$$\begin{aligned} ev : \mathcal{M} &\rightarrow \Lambda \\ u &\mapsto u(0) \end{aligned}$$

(where $u(s)$ stands for $t \in S^1 \mapsto u(s, t) \in M$) and the right action of \mathbb{R} defined by

$$(u \cdot \sigma)(s, t) = u(s + \sigma, t).$$

One checks easily that $u \in \mathcal{M}$ implies $u \cdot \sigma \in \mathcal{M}$.

Proposition 4.8. *For $u \in \mathcal{M}$, any cluster point of $u(s)$ in Λ when $s \rightarrow \pm\infty$ is a critical point of the action functional.*

Proof. This is a little trickier than in finite dimension because we do not have a flow on $ev(\mathcal{M})$ (no Cauchy-Lipschitz theorem). Let $s_k \rightarrow +\infty$ be a sequence such that $u(s_k) \rightarrow \lambda \in \Lambda$. Certainly λ belongs to $ev(\mathcal{M})$. Indeed, since \mathcal{M} is compact and taking a subsequence if necessary, $u \cdot s_k \rightarrow v \in \mathcal{M}$ and $v(0) = \lambda$. The convergence being uniform on each compact interval $[-T, +T]$ and according to proposition 4.1 1), we deduce :

$$\sup(\Phi_H \circ v)(\mathbb{R}) = \sup_{T \rightarrow +\infty} \lim_k \Phi_H(u(s_k - T)) = \lim_{s \rightarrow +\infty} \Phi_H(u(s)),$$

$$\inf(\Phi_H \circ v)(\mathbb{R}) = \inf_{T \rightarrow +\infty} \lim_k \Phi_H(u(s_k + T)) = \lim_{s \rightarrow +\infty} \Phi_H(u(s)).$$

Then $(\Phi_H \circ v)(\mathbb{R})$ contains a single point $\Phi_H(\lambda)$. Since v is a gradient line, this implies that λ is a critical point. \square

Corollary 4.9.

- 1) If all 1-periodic solutions of X_{H_t} are non-degenerate, then for any $u \in \mathcal{M}$, $u(s) \rightarrow u_{\pm}$ in Λ when $s \rightarrow \pm\infty$.
- 2) $\frac{\partial u}{\partial s}(s, t)$ tends to 0 when $s \rightarrow \pm\infty$ uniformly in $t \in S^1$.

Proof.

1) We consider only $s \rightarrow +\infty$. We already know that there are cluster points (proposition 4.1). The existence of a limit reduces to the uniqueness of cluster points when $s \rightarrow +\infty$. This set of cluster points is

$$A = \bigcap_n \left(\overline{\{u(s) \mid s \geq n\}} \right)$$

and, as $ev(\mathcal{M})$ is compact, it is a decreasing intersection of compact connected sets. Therefore it is connected. But, by the previous proposition, it is included in the finite set $\text{crit}(\Phi_H)$. So A contains only one point.

2) For $u \in \mathcal{M}$ we have:

$$-\frac{\partial u}{\partial s}(s, t) = \left[J(u) \frac{\partial u}{\partial t} + \nabla H(t, u) \right] (s, t)$$

and, according to 1), it tends to

$$[J(x_{\pm})\dot{x}_{\pm} + \nabla H(t, x_{\pm})](t)$$

uniformly in t . But this limit is 0 since x_{\pm} is a 1-periodic solution of the Hamilton equation. \square

In other words any $u \in \mathcal{M}$ belongs to some $\mathcal{M}(x_-, x_+)$, the set of gradient lines of Φ_H descending from x_- to x_+ . The right action of \mathbb{R} keeps $\mathcal{M}(x_-, x_+)$ invariant and an equivalence class \bar{u} is called a *connecting orbit*. Now it is possible to state and prove an analogue of theorem 3.5 2).

Theorem 4.10. (Broken connecting orbit theorem).

Let $u_n \rightarrow u$ be a converging sequence in \mathcal{M} , with $u_n \in \mathcal{M}(x_-, x_+)$ for all n . Then there exist a subsequence (still denoted by u_n) and:

- 1-periodic solutions of X_{H_t} , $x_0 = x_-, x_1, \dots, x_{\ell}, x_{\ell+1} = x_+$
- for each $k \in \{0, \dots, \ell\}$, a sequence of real numbers $(s_n^k)_n$ and a L^2 -gradient line $u^k \in \mathcal{M}(x_k, x_{k+1})$

such that $u_n \cdot s_n^k \rightarrow u^k$ when $n \rightarrow \infty$.

The $(\ell + 1)$ -uple of connecting orbits $(\bar{u}^0, \dots, \bar{u}^{\ell})$ is called a *broken connecting orbit*. The proof goes as in theorem 3.5 where the Morse function is replaced by the action functional. Intentionally we avoided making use of Morse models. We only used that the set of gradient lines is compact and the finiteness of the critical set, which is true in our infinite dimensional setting.

5 Floer Homology

Here is an “easy” case of Floer’s theorem we are aiming at.

Theorem 5.1. *Let (M, ω) be a closed symplectic manifold. Let $H : \mathbb{R} \times M \rightarrow \mathbb{R}$ be a Hamiltonian 1-periodic in time and X_H be its associated vector field. We assume the following:*

- 1) $\langle [\omega], \pi_2(M) \rangle = 0$,
- 2) $\langle c_1(M), \pi_2(M) \rangle = 0$,
- 3) the 1-periodic solutions of X_{H_t} are non-degenerate.

Then the number of 1-periodic solutions of the Hamilton equation is at least $\dim_{\mathbb{Z}/2} H_(M, \mathbb{Z}/2)$.*

Condition 2) means that the tangent fiber space, equipped with any complex structure J which is ω -compatible, is trivial as a complex fiber bundle over any 2-sphere in M and this condition does not depend on the chosen J . Actually there are more precise conclusions like Morse inequalities. But they require defining an index for each 1-periodic solution of the Hamilton equation. One could think of 1-periodic solutions as critical points of the action functional and look at its Morse index. Unfortunately this does not work. When $M = \mathbb{R}^{2n}$, an easy calculation by Fourier transform shows that both the index and co-index are infinite. The right index to consider is the *Conley-Zehnder index*.

5.1 The Conley-Zehnder index (case of \mathbb{R}^{2n})

Let $\dot{x} = J\nabla H(t, x)$ be the Hamiltonian differential equation under consideration and let $(x(t))$ be a non-degenerate 1-periodic solution. The linearized differential equation along $(x(t))$ reads

$$\dot{\xi} = JS(t)\xi$$

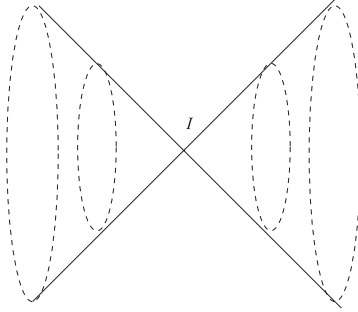
where $S(t) := D^2H(t, x(t))$, a symmetric matrix 1-periodic in t . Let (R_t) be its linear flow. It is a path in $\mathrm{Sp}(2n)$, the $2n \times 2n$ symplectic matrices, such that $R_0 = I$ and $\dot{R}_t = JS(t)R_t$. The non-degeneracy condition of the 1-periodic $(x(t))$ reads :

$$R_1 \in \mathrm{Sp}(2n)^* := \{A \in \mathrm{Sp}(2n) \mid \det(I - A) \neq 0\}.$$

The *Maslov cycle* is

$$\Sigma := \{A \in \mathrm{Sp}(2n) \mid \det(I - A) = 0\}.$$

It is a hypersurface whose singularities have codimension ≥ 3 in $\mathrm{Sp}(2n)$ (compare [2]).

Figure 7: The Maslov cycle in $\text{Sp}(2)$ near I .

Proposition 5.2. [15][45]

1) $\text{Sp}(2n)^*$ has two connected components:

$$\text{Sp}(2n)^\pm = \{A \mid \det(I - A) \in \mathbb{R}^\pm\}.$$

2) Any loop in $\text{Sp}(2n)^*$ is contractible in $\text{Sp}(2n)$.

Another proof is given in the appendix.

Theorem 5.3. (Conley-Zehnder) [15][45]

1) Given a path (R_t) in $\text{Sp}(2n)$ from I to $R_1 \in \text{Sp}(2n)^*$, there exists a “natural” index $\mu((R_t)) \in \mathbb{Z}$ which is invariant under homotopy as long as R_1 remains in $\text{Sp}(2n)^*$. Moreover it classifies the homotopy classes of such paths (R_t) .

2) We have the following sign formula

$$\text{sign } \det(I - R_1) = (-1)^{n - \mu((R_t))}.$$

3) When S is symmetric, independent of t with $\|S\| < 2\pi$, then

$$\mu((\exp tJS)) = -\text{ind}(S) + n,$$

where $\text{ind}(S)$ denotes the number of negative eigenvalues of S .

This index is “natural” in the sense that it is invariant under conjugacy and additive with respect to direct sums. In [6] it is related to a *symplectic translation number*.

Sketch of proof.

1) The polar decomposition yields a homeomorphism

$$\mathrm{Sp}(2n) \rightarrow Q \times \mathrm{U}(n)$$

where $\mathrm{U}(n)$ is the unitary group, the intersection $\mathrm{Sp}(2n) \cap \mathrm{O}(2n)$, and where Q is the subgroup of symmetric, positive definite, symplectic matrices. The space Q is contractible and

$$\det_{\mathbb{C}} : \mathrm{U}(n) \rightarrow S^1$$

induces an isomorphism on π_1 . Here is a hint for the contractibility of Q : for $A \in Q$, let E_{λ} be the eigenspace of eigenvalue $\lambda \geq 1$ and set $F_{\lambda} := E_{\lambda} \oplus E_{\lambda^{-1}}$; F_{λ} is a symplectic J -invariant subspace; then A reads in an unique way $A = e^{JS}$ where S is symmetric and anticommutes with J . Thus a contraction is provided by $(A, t) \mapsto e^{tJS}$, $t \in [0, 1]$.

From the polar decomposition we get a continuous map

$$\rho : \mathrm{Sp}(2n) \rightarrow S^1,$$

which is the composition of the projection $\mathrm{Sp}(2n) \rightarrow \mathrm{U}(n)$ followed by $\det_{\mathbb{C}}$. It induces an isomorphism on π_1 . In both connected components $\mathrm{Sp}(2n)^{\pm}$ we choose base points:

$$\begin{aligned} W_+ &:= \mathrm{diag}(-1, \dots, -1) \in \mathrm{Sp}(2n)^+ \\ W_- &:= \mathrm{diag}(2, -1, \dots, \frac{1}{2}, -1, \dots) \in \mathrm{Sp}(2n)^- \end{aligned}$$

We have $\rho(W_+) = (-1)^n$ and $\rho(W_-) = (-1)^{n-1}$. Given a path (R_t) , one connects R_1 to W_{\pm} by a path $\alpha \in \mathrm{Sp}(2n)^*$ yielding a new path (R'_t) from I to W_{\pm} in the same homotopy class. In lifting $t \mapsto \rho(R'_t)$ to the universal cover of S^1 we get $\rho(t) = e^{ia(t)}$ with $a(0) = 0$ and $a(1) \in \pi\mathbb{Z}$. We now can define:

$$\mu((R_t)) = \frac{a(1)}{\pi}.$$

This integer does not depend on the choice of α because a change of α has the effect of composing (R'_t) with a loop λ in $\mathrm{Sp}(2n)^*$ and, due to proposition 5.2 2), it does not change $\mu((R_t))$. Therefore this integer is an invariant of the homotopy class of (R_t) as long as R_1 remains in $\mathrm{Sp}(2n)^*$. The fact that it is a complete invariant of the homotopy class is a direct consequence of ρ being a π_1 -isomorphism. Therefore we have the following invariance by conjugation

$$\mu((P_t R_t P_t^{-1})) = \mu((R_t))$$

where $t \mapsto P_t$ is any path in $\mathrm{Sp}(2n)$.

2) For the sign formula we can take $R_1 = W_{\pm}$ and any path α from I to W_{\pm} . Indeed composing α by a loop based at W_{\pm} changes $\mu((R_t))$ by an even integer. Then calculating the invariant associated with a simple path from I to

W_{\pm} yields the sign formula.

3) We only check this point on examples. For instance when $n = 1$ and $S = \begin{pmatrix} \varepsilon & 0 \\ 0 & \varepsilon \end{pmatrix}$, we have $\exp JS = \begin{pmatrix} \cos \varepsilon & -\sin \varepsilon \\ \sin \varepsilon & \cos \varepsilon \end{pmatrix}$ and $\rho(\exp JS) = e^{i\varepsilon}$. If $2\pi > \varepsilon > 0$, then, with the above notation, $R'_1 = W_+$ and $a(1) = +\pi$. If $-2\pi < \varepsilon < 0$, then $a(1) = -\pi$. In both cases we have

$$\mu((\exp tJS)) = -\text{ind}(S) + n.$$

When $S = \begin{pmatrix} -\varepsilon & 0 \\ 0 & +\varepsilon \end{pmatrix}$, JS has two real eigenspaces and $\mu((\exp tJS)) = 0$. The index formula is still true in that case. \square

Remark 5.4. If (U_t) is a loop in $U(n)$ based at I with $\det(U_t) = e^{ia(t)}$ and $a(0) = 0$, then $a(1)/\pi$ is the usual Maslov index of this loop and it belongs to $2\mathbb{Z}$.

Definition 5.5. In the same setting as above the Conley-Zehnder index is defined by

$$i_{CZ}((R_t)) = -\mu((R_t)).$$

The μ -index is closely related to the Maslov index while the i_{CZ} -index is closer to the Morse index.

5.2 Relation with a Fredholm index

Let us consider

$$F := \frac{\partial}{\partial s} + J \frac{\partial}{\partial t} + S(s, t)$$

where $S : \mathbb{R} \times S^1 \rightarrow \text{End}(\mathbb{R}^{2n})$ is continuous. When $s \rightarrow \pm\infty$, $S(s, t)$ is assumed to tend to a symmetric matrix $S_{\pm}(t)$. Let R_t^{\pm} be the solution of $\dot{R} = JS_{\pm}(t) \cdot R$ with initial condition $R_0^{\pm} = I$.

Theorem 5.6. *If $\det(I - R_1^{\pm}) \neq 0$, then $F : W^{1,2}(\mathbb{R} \times S^1, \mathbb{R}^{2n}) \rightarrow L^2(\mathbb{R} \times S^1, \mathbb{R}^{2n})$ is a Fredholm operator. Its Fredholm index is given by*

$$\begin{aligned} \text{ind}(F) &= \mu((R_1^+)) - \mu((R_1^-)) \\ &= i_{CZ}((R_1^-)) - i_{CZ}((R_1^+)). \end{aligned}$$

We recall that $W^{1,2}$ is the Sobolev space of functions in L^2 whose weak derivatives are in L^2 . It is a Hilbert space for a suitable scalar product. The operator

F is Fredholm if it has a closed range and if its kernel and cokernel have a finite dimension. By definition its Fredholm index is

$$\text{ind}(F) = \dim \ker F - \dim \text{coker } F.$$

Proof. The Fredholm property and the index do not change by adding a compact operator, for instance in modifying S among the bounded matrix functions. The right-hand side of the index formula does not change as long as R_1^\pm remains in $\text{Sp}^*(2n)$. In this way we are reduced to the case when $S(s, t)$ is symmetric and independent of t . For every $s \in \mathbb{R}$,

$$A(s) := J \frac{\partial}{\partial t} + S(s) : W^{1,2}(S^1, \mathbb{R}^{2n}) \rightarrow L^2(S^1, \mathbb{R}^{2n})$$

is self-adjoint. Its spectrum is discrete. It contains 0 if and only if there exists some non-trivial 1-periodic solution of

$$\frac{\partial \xi}{\partial t} = JS(s)\xi$$

that is $\det(I - R_{s1}) = 0$, where R_{st} is the solution of $\dot{R} = JS(s)R$, $R_{s0} = I$. It is known ([8] Section 17, [43] Thm. A) that the Fredholm index of F is the *spectral flow* of the family $(A(s))$ when s goes from $-\infty$ to $+\infty$, that is the algebraic number of eigenvalues crossing 0. Finally $\text{ind}(F)$ is the algebraic number of crossings of $(R_{s1})_{-\infty}^{+\infty}$ with the Maslov cycle when it is endowed with a suitable co-orientation (cf. Appendix). It is “clear” from Figure 8 that this number is the desired one $\mu((R_t^+)) - \mu((R_t^-))$.

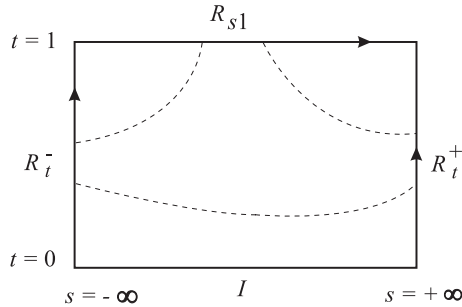


Figure 8: The dot lines are $\{(s, t) \mid \det(I - R_{s,t}) = 0, t > 0\}$.

Of course one needs to be careful along $t = 0$. But the index of F depends only on the path (R_{s1}) and the Conley-Zehnder indices of the vertical sides

do not change if their endpoints are kept fixed. This remark allows us to take $S(s, t)$ invertible for all s and small t . In that case $\{t = 0\}$ becomes isolated in $\{\det(I - R_{st}) = 0\}$ and Figure 8 tells us the truth. \square

5.3 The Conley-Zehnder index of a 1-periodic orbit

Let us consider (M, ω, H) as in theorem 5.1. It is possible to define $i_{CZ}(x)$ when x is a non-degenerate 1-periodic orbit of the Hamilton equation. We choose an almost complex structure J which is ω -compatible. The complex fiber bundle (x^*TM, J) is trivial and, in fact, this trivialization is unique up to homotopy. Indeed, according to 2) from theorem 5.1, it does not depend on the contraction chosen for x . Having trivialized (x^*TM, J) , the linearized differential equation reads as in the \mathbb{R}^{2n} - case and $i_{CZ}(x)$ is defined as before.

5.4 Application to the gradient lines of the action

Notations. For a 1-periodic Hamiltonian H , P_H will denote the set of contractible 1-periodic solutions of the Hamilton equation $\dot{x} = X_H(t, x)$. By assumption they are all non-degenerate. We have

$$P_H = \text{crit } \Phi_H$$

where Φ_H is the action functional defined on the smooth contractible loops. For $x_-, x_+ \in P_H$, $\mathcal{M}(x_-, x_+)$ will denote the set of gradient lines descending from x_- to x_+ . That means:

$$\left\{ \begin{array}{l} u \in \mathcal{M}(x_-, x_+) \iff \\ u : \mathbb{R} \times S^1 \rightarrow T^{2n} \text{ is smooth} \\ \frac{\partial u}{\partial s} + J(u) \frac{\partial u}{\partial t} + \nabla H(t, u) = 0 \\ u \rightarrow x_{\pm} \text{ if } s \rightarrow \pm\infty. \end{array} \right.$$

Remark 5.7. About the analytic setting.

The general aim is to describe the set $\mathcal{M}(x_-, x_+)$ as a submanifold of some Banach manifold on which the *perturbed Cauchy-Riemann equation* would be a regular Fredholm equation. This can be achieved when working with $W^{1,p}$, $p > 2$, the space of functions in L^p whose first derivatives are in L^p . When the source is 2-dimensional, such a function is continuous (it can be non-continuous if $p = 2$). As we do not intend to deal with this sort of difficulty, in the sequel we shall restrict ourselves to the case where M is the $2n$ -torus as a differentiable manifold; as it is said below, in that case it is sufficient to work with $W^{1,2}$. But except for the Sobolev setting we will argue as in the general case.

Indeed, when $M = T^{2n} = \mathbb{R}^{2n}/\mathbb{Z}^{2n}$ and if $\pi : \mathbb{R}^{2n} \rightarrow T^{2n}$ denotes the projection, it is easy to define $W^{1,2}(\mathbb{R} \times S^1, M, x_-, x_+)$; it is a set of u 's of the form $u = \pi \circ \tilde{u}$, where \tilde{u} belongs to $W_{loc}^{1,2}(\mathbb{R} \times S^1, \mathbb{R}^{2n})$ with some limit conditions at infinity we are explaining now. First we recall for $\tilde{u} \in W_{loc}^{1,2}(\mathbb{R} \times S^1, \mathbb{R}^{2n})$, according to the Fubini theorem, for almost all $s \in \mathbb{R}$ the restriction $\tilde{u}(s) = \tilde{u}|_{\{s\} \times S^1}$ lies in $W^{1,2}(S^1, \mathbb{R}^{2n})$, which is a space of continuous functions; the corresponding $u(s) : S^1 \rightarrow T^{2n}$ is continuous accordingly. This enables us to define $W^{1,2}(\mathbb{R} \times S^1, T^{2n}, x_-, x_+)$ by the following requirements: there is some lifting \tilde{x}_\pm of x_\pm to \mathbb{R}^{2n} such that $\tilde{u} - \tilde{x}_\pm$ belongs to $W^{1,2}(\mathbb{R}_\pm \times S^1, \mathbb{R}^{2n})$. It makes sense because the loop x_\pm is homotopic to zero.

As a consequence, for almost every $t \in S^1$, $\tilde{u}(s, t) - \tilde{x}_\pm(t)$ is continuous and tends to 0 when s goes to $\pm\infty$ ([10] corollaire VIII.8). Therefore the lifting x_\pm is uniquely determined by the choice of \tilde{u} representing u ; in particular the choice of \tilde{u} is unique up to covering transformation. Now it is clear that $W^{1,2}(\mathbb{R} \times S^1, T^{2n}, x_-, x_+)$ is a Banach manifold modelled on $W^{1,2}(\mathbb{R} \times S^1, \mathbb{R}^{2n})$.

Proposition 5.8. $\mathcal{M}(x_-, x_+)$ is defined by a Fredholm equation in the sense that, at every $u \in \mathcal{M}(x_-, x_+)$, the linearized equation $F_u \xi = 0$ is Fredholm. Moreover the index of F_u is

$$\text{ind}(F_u) = i_{CZ}(x_-) - i_{CZ}(x_+).$$

Proof. For $u \in \mathcal{M}(x_-, x_+)$ and ξ , a vector field along u tangent to T^{2n} ,

$$F_u \xi := D_s \xi + J(u) D_t \xi + (D_\xi J(u)) \cdot \frac{\partial u}{\partial t} + D_\xi(\nabla H(t, u))$$

where D_s and D_t stand for the covariant derivative with respect to $\frac{\partial u}{\partial s}$ and $\frac{\partial u}{\partial t}$ respectively. Let $R = (Z_1, \dots, Z_{2n})$ be an orthonormal J -invariant framing of $u^*(T(T^{2n}))$ (that is $JR = R$ up to a permutation) which extends over the limit orbits x_+ and x_- . If $\xi = R\zeta$, with $\zeta(s, t) \in \mathbb{R}^{2n}$, we get

$$\left(R^* F_u \right) \zeta = \frac{\partial \zeta}{\partial s} + J \frac{\partial \zeta}{\partial t} + L \cdot \zeta$$

where $L \cdot \zeta$ is an operator of order 0, that is $L(s, t) \in \text{End}(\mathbb{R}^{2n})$.

Claim. When s goes to $\pm\infty$, the limit of $L(s, t)$ is a symmetric matrix.

Proof of the claim. The (i, j) -entry of $L(s, t)$ is given by the value at (s, t) of the following sum:

$$\langle Z_i, D_s Z_j + J D_t Z_j + D_{Z_j}(J(u)) \cdot \frac{\partial u}{\partial t} + D_{Z_j} \nabla H \rangle.$$

According to corollary 4.9 the limit when $s \rightarrow \pm\infty$ is

$$\ell_{ij}^\pm(t) := \langle Z_i, JD_t Z_j + D_{Z_j}(J) \cdot \dot{x}_\pm + D_{Z_j} \nabla H \rangle |_{x_\pm(t)} .$$

From $\dot{x}_\pm = J(x_\pm) \nabla H(t, x_\pm)$ and $J.D_{Z_j}(J) = -D_{Z_j}(J).J$ we derive :

$$\begin{aligned} \ell_{ij}^\pm(t) &= \langle Z_i, J(D_X Z_j - D_{Z_j} X) \rangle \\ &= \langle Z_i, J[X, Z_j] \rangle = -\omega(Z_i, [X, Z_j]), \end{aligned}$$

where $X := J \nabla H$. Since X is a Hamiltonian vector field its flow preserves ω . Then

$$\begin{aligned} 0 &= X.\omega(Z_i, Z_j) \\ &= (L_X \omega)(Z_i, Z_j) + \omega([X, Z_i], Z_j) + \omega(Z_i, [X, Z_j]) \\ &= \omega(Z_i, [X, Z_j]) - \omega(Z_j, [X, Z_i]), \end{aligned}$$

hence L is symmetric. \square

As a consequence the linearized equation $(R^* F_u) \zeta = 0$ takes the form required in 5.2. From theorem 5.6 we deduce that

$$R^* F_u : W^{1,2}(\mathbb{R} \times S^1, \mathbb{R}^{2n}) \rightarrow L^2(\mathbb{R} \times S^1, \mathbb{R}^{2n})$$

is a Fredholm operator and its index is $i_{CZ}(x_-) - i_{CZ}(x_+)$. \square

In general the Fredholm property is not sufficient in order to deduce that $\mathcal{M}(x_-, x_+)$ is a finite dimensional manifold. But we observe that H and J play the role of parameters in the equation of $\mathcal{M}(x_-, x_+)$.

Definition 5.9. *(H, J) is a regular pair if, for every $u \in \mathcal{M}(x_-, x_+)$, F_u is surjective.*

It is easy to prove that there are enough possible deformations of (H, J) that for a pair to be regular is a generic property in Baire's sense. This fact follows from the infinite dimensional version of Sard's theorem due to Smale, which applies to (C^∞) - nonlinear Fredholm maps between Banach spaces ([51]).

Remark 5.10. If we start with a Hamiltonian function H whose 1-periodic solutions are non-degenerate and with an arbitrary ω -compatible almost complex structure J , then approximating (H, J) by a regular pair (\tilde{H}, \tilde{J}) close enough to (H, J) does not change the number of 1-periodic solutions nor their indices since they can be followed along the approximation, according to the implicit function theorem (compare subsection 4.3). Moreover I think it is possible to take $\tilde{H} = H$, but after the previous remark this question is irrelevant for theorem 5.1.

Corollary 5.11. *When (H, J) is regular, $\mathcal{M}(x_-, x_+)$ is a finite dimensional manifold, whose dimension is $i_{CZ}(x_-) - i_{CZ}(x_+)$.*

Proof. Apply the implicit function theorem. \square

As in Morse theory we have a natural right action of \mathbb{R} onto the manifold $\mathcal{M}(x_-, x_+)$ by translation of the variable s . In particular when $i_{CZ}(x_-) - i_{CZ}(x_+) = 0$, $\mathcal{M}(x_-, x_+)$ is empty since if it contains a point it would contain a line.

Theorem 5.12. *There exists a natural compactification of $\mathcal{M}(x_-, x_+)/\mathbb{R}$ made of adding the broken L^2 -gradient lines from x_- to x_+ . It is a manifold (with boundary and corners) of dimension $i_{CZ}(x_-) - i_{CZ}(x_+) - 1$.*

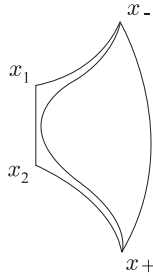


Figure 9: A broken L^2 -gradient line with two vertices.

In the example of figure 9 we have $i_{CZ}(x_-) > i_{CZ}(x_1) > i_{CZ}(x_2) > i_{CZ}(x_+)$. Sometimes this theorem is called the *glueing theorem* and it corresponds to a piece of work in Floer's theory.

Corollary 5.13. *If $i_{CZ}(x_-) = i_{CZ}(x_+) + 1$, $\mathcal{M}(x_-, x_+)/\mathbb{R}$ is a compact 0-dimensional manifold and hence a finite set whose cardinality mod 2 is denoted by $n(x_-, x_+)$.*

5.5 The Floer complex

Given a regular pair (H, J) , the vector spaces of the Floer complex are:

$$CF_k(H, J) = \bigoplus_{x \in P(H), i_{CZ}(x) = k} (\mathbb{Z}/2\mathbb{Z}) \cdot x$$

The boundary operator $\partial_k : CF_k(H, J) \rightarrow CF_{k-1}(H, J)$ is defined by:

$$\partial_k(x) = \sum_{i_{CZ}(y) = k-1} n(x, y) y.$$

Theorem 5.14.

$$\partial_k \circ \partial_{k+1} = 0.$$

Using theorem 5.12, the proof of theorem 5.14 goes as in Morse theory. This allows us to define homology groups

$$HF_k(H, J) = \frac{\ker \partial_k}{\text{im } \partial_{k+1}}.$$

The following theorem claims that the Floer homology does not depend on the chosen regular pair (H, J) . It is similar to theorem 3.10.

Theorem 5.15. *If (H_0, J_0) and (H_1, J_1) are two regular pairs then there exists a canonical isomorphism*

$$HF_*(H_0, J_0) \rightarrow HF_*(H_1, J_1).$$

The last result we will mention and which implies theorem 5.1 is the following.

Theorem 5.16. *Let H be a Morse function independent of time and J an almost complex structure ω -compatible. We assume the Smale transversality condition: for the Riemannian metric $\langle \xi, \eta \rangle = \omega(\xi, J\eta)$ the stable and unstable manifolds are transversal in M .*

Then, for $\tau \in]0, 1]$ small enough, the following holds:

- 1) $(\tau H, J)$ is regular.
- 2) $HF_*(\tau H, J) \cong H_{*+n}(M, \mathbb{Z}/2\mathbb{Z})$.

Proof. We can start with H so C^2 -small that lemma 2.2 applies. Then the only 1-periodic orbits of $\dot{x} = J\nabla H(x)$ are the critical points of H . In that case the linearized Hamilton equation at a critical point x_0 is

$$\dot{\xi} = JD^2H(x_0)\xi.$$

If $\|D^2H(x_0)\| < 2\pi$ the Conley-Zehnder index is $i_{CZ}(x_0) = \text{ind}(x_0) - n$. Let x_- and x_+ be two critical points of H with $i_{CZ}(x_-) = i_{CZ}(x_+) + 1$. We have to find all the solutions of

$$\frac{\partial u}{\partial s} + J(u) \frac{\partial u}{\partial t} + \nabla H(u) = 0$$

such that $u \rightarrow x_{\pm}$ when s goes to $\pm\infty$. For this purpose we follow Salamon-Zehnder ([45], p. 1342).

Lemma 5.17. *There exists $\tau_0 > 0$ such that, for every $\tau \leq \tau_0$, any $u \in \mathcal{M}(x_-, x_+)$, which is τ -periodic in t , is independent of t .*

Proof. If this is not true, there exists a sequence $u_k \in \mathcal{M}(x_-, x_+)$ where u_k is periodic in t with a minimal period τ_k such that $0 < \tau_k \leq 1/k$. According to the broken connecting orbit theorem 4.10, taking a subsequence and a suitable shift of s if necessary, u_k converges to a piece $u \in \mathcal{M}$ of a broken connecting orbit γ from x_- to x_+ . Obviously u (and all pieces of γ as well) is 0-periodic in t and hence it is a solution of

$$\frac{\partial u}{\partial s} + \nabla H(u) = 0 \quad (1)$$

that is a usual gradient line of H in M . Under the Smale transversality condition, there are no usual broken connecting orbits from x_- to x_+ in M which are really broken, because their Morse indices differ by 1. As a consequence $u_k \rightarrow u \in \mathcal{M}(x_-, x_+)$ which is a usual gradient line. For such an u , there are two linearized operators. The first one F_u^0 is obtained in linearizing equation (1) of gradient lines of H . The second one F_u^τ , depending on a parameter τ , is obtained in linearizing the Cauchy-Riemann equation

$$\frac{\partial u}{\partial s} + J(u) \frac{\partial u}{\partial t} + \nabla H(u) = 0 \quad (2)$$

where u is τ -periodic in t , that is $u : \mathbb{R} \times \mathbb{R}/\tau\mathbb{Z} \rightarrow M$ and $x_\pm = \lim_{s \rightarrow \pm\infty} u(s, t)$.

Let us be more precise. As in proposition 5.8 we use a framing R of u^*TM . Because u does not depend on t , one can choose R so that it does not either. The equation $F_u^\tau \xi = 0$ becomes

$$(R^* F_u^\tau) \zeta := \frac{\partial \zeta}{\partial s} + J \frac{\partial \zeta}{\partial t} + S(s) \zeta = 0$$

for $\zeta : \mathbb{R} \times \mathbb{R}/\tau\mathbb{Z} \rightarrow \mathbb{R}^{2n}$. The equation $F_u^0 \xi = 0$ becomes

$$(R^* F_u^0) \zeta := \frac{d\zeta}{ds} + S(s) \zeta = 0$$

for $\zeta : \mathbb{R} \rightarrow \mathbb{R}^{2n}$. Here $S(s)$ is an endomorphism of \mathbb{R}^{2n} having a limit at infinity (indeed a symmetric endomorphism).

Sublemma.

1) $F_u^0 : W^{1,2}(\mathbb{R}, u^*TM) \rightarrow L^2(\mathbb{R}, u^*TM)$ is surjective and its Fredholm index is $\text{ind}(x_-) - \text{ind}(x_+)$.

If τ is small enough, the two following conclusions hold:

2) $\ker F_u^0 = \ker F_u^\tau$.

3) F_u^τ is surjective.

References for proofs

1) This is proven with all details in [44], thm. 3.3.

2) For $\zeta : \mathbb{R} \times \mathbb{R}/\tau\mathbb{Z} \rightarrow \mathbb{R}^{2n}$, let $\mu(\zeta)$ be its mean with respect to the t -variable. On the one hand, $(R^*F_u^\tau)\zeta = 0$ yields $(R^*F_u^0)\mu(\zeta) = 0$. On the other hand if $\mu(\zeta) = 0$, there is an elementary estimate

$$\|\zeta\|_{L^2} \leq \tau \left\| \frac{\partial \zeta}{\partial t} \right\|_{L^2}$$

from which one can derive :

$$\left\| \frac{\partial \zeta}{\partial t} \right\|_{L^2} \leq \|(R^*F_u^\tau)\zeta\|_{L^2} + c\tau \left\| \frac{\partial \zeta}{\partial t} \right\|_{L^2}$$

where c is a uniform bound of $\|S(s)\|$ (compare [45], p. 1327-1328). When τ is small enough, we have $c\tau < 1$. In that case $(R^*F_u^\tau)\zeta = 0$ and $\mu(\zeta) = 0$ imply $\zeta = 0$; or: $(R^*F_u^\tau)\zeta = 0$ implies $\zeta = \mu(\zeta)$, which is independent of t . Finally $\xi \in \ker F_u^\tau$ implies $\xi \in \ker F_u^0$. The converse being obvious we have the wanted conclusion.

3) According to 1) we have

$$\begin{aligned} \dim \ker F_u^0 &= \text{ind}(x_-) - \text{ind}(x_+) \\ &= i_{CZ}(x_-) - i_{CZ}(x_+) \\ &= \text{ind}(F_u^\tau). \end{aligned}$$

According to 2) it is also the dimension of $\ker F_u^\tau$. Thus F_u^τ is surjective. \square

We return to the proof of lemma 5.17. Let us fix some τ such that points 2) and 3) in the above sublemma are valid. Since the Fredholm index of F_u^τ is 1, by the implicit function theorem, u is isolated in $\mathcal{M}_\tau(x_-, x_+)$ up to a shift in s , where $\mathcal{M}_\tau(x_-, x_+)$ is the space of solutions of the perturbed Cauchy Riemann equation, which are τ -periodic in time (it makes sense even if τ is not a divisor of 1, because the Hamiltonian is autonomous) and whose limits at infinity are respectively x_- and x_+ . Moreover u is uniformly isolated in $\mathcal{M}_{\tau'}(x_-, x_+)$ for all τ' close enough to τ , let us say $\tau' \in]\tau - \varepsilon, \tau + \varepsilon[$; for that we use the implicit function theorem with a parameter. When k is large, the period τ_k is small, some integer multiple τ'_k of τ_k lies in $]\tau - \varepsilon, \tau + \varepsilon[$ and u_k is close to u . Hence u is not isolated in $\mathcal{M}_{\tau'_k}(x_-, x_+)$, leading to a contradiction. So the lemma is proven. \square

We now finish the proof of theorem 5.16. Let us consider the Hamiltonian τH and $u_\tau(s, t) := u(\tau s, \tau t)$. One checks easily that:

1. u_τ is 1-periodic in t iff u is τ -periodic,
2. u_τ is a solution of the Cauchy-Riemann equation perturbed by τH with x_\pm as boundary condition iff $u \in \mathcal{M}(x_-, x_+)$.
3. F_{u_τ} is surjective iff F_u^τ is surjective.

Let u_τ be a 1-periodic solution of

$$\frac{\partial u}{\partial s} + J(u) \frac{\partial u}{\partial t} + \nabla \tau H(u) = 0$$

with x_\pm as boundary conditions. Let u be the associated τ -periodic map as above. According to the sublemma, if τ is small enough the pair $(\tau H, J)$ is regular.

In choosing τ_0 so that lemma 5.17 is valid, we deduce that u is 0-periodic, hence u_τ also is. In other words the only 1-periodic solutions of

$$\frac{\partial u}{\partial s} + J(u) \frac{\partial u}{\partial t} + \nabla \tau H(u) = 0$$

with x_\pm as boundary conditions are the solutions of

$$\frac{\partial u}{\partial s} + \nabla \tau H(u) = 0$$

that is the usual gradient lines from x_- to x_+ . This amounts to saying that the Floer complex $CF_*(\tau H, J)$ coincides with the Morse complex up to a translation of the grading. \square

Corollary 5.18. *Let us consider the same data and assumptions as in theorem 5.1. Let us denote by p_k the number of 1-periodic orbits with Conley-Zehnder index k and by b_k the k^{th} Betti number over $\mathbb{Z}/2\mathbb{Z}$. The Morse inequalities are valid:*

$$p_k - p_{k-1} + \dots \geq b_{n+k} - b_{n+k-1} + \dots + (-1)^{k+n} b_0.$$

Final comments

There are many other problems where a sort of Floer homology or cohomology is involved; one speaks of cohomology when a product is defined and used. Such groups, even when there are well defined, are not always computable. It is the same situation as in classical algebraic topology; as there, even if HF_* or HF^* is not computable, some deep results can be deduced for instance in arguing with

exact sequences. Here are a few examples where various Floer (co)-homologies appear.

1) Fixed points of symplectic diffeomorphisms which are not Hamiltonian (see [16]). In [47] P. Seidel succeeded in computing the Floer homology of a Dehn twist along a multi-curve on a surface.

2) Intersections of Lagrangian submanifolds [21]. Here we have to recall that this problem (and paper) is the historical one because it is the first which has been solved by means of Floer homology. In fact the problem of fixed points of symplectic diffeomorphisms could be regarded as a particular case of this one: take the graph of the diffeomorphism under consideration and look at its intersection with the diagonal. Due to the topological restrictions in [21], Floer found more efficient to create another Floer homology for the fixed points problem in [24].

3) Picard-Lefschetz theory. Of particular interest is the generalization of the 2-dimensional Dehn twist discovered by V.I. Arnold in high dimension, namely the *Dehn twist along a Lagrangian sphere* [5]. It appears as the monodromy diffeomorphism of a complex Morse singularity. P. Seidel formulated a symplectic analogue of the holomorphic Morse functions and revealed a long exact sequence involving the Floer homology of a pair of exact Lagrangian submanifolds and the Dehn twist along an exact Lagrangian sphere [48].

4) The Weinstein conjecture [57]: Any closed hypersurface Σ of *contact type* in a symplectic manifold (M, ω) carries a closed *characteristic (contact type)* means that along Σ there is a transversal vector field dilating ω . In complete generality it is still an open question. The *linear case* $(M, \omega) = (\mathbb{R}^{2n}, \omega_0)$ was proven by C. Viterbo [55]; the proof was based on a linking phenomenon in some loop space. Since this time many other cases have been solved by means of various tools including pseudoholomorphic techniques. In [56] Viterbo succeeded in putting the Weinstein conjecture into the Floer homology world. It allowed him to recover the known cases and to solve new ones like: $M = T^*N$ where N is 1-connected, or M is a *sub-critical* Stein manifold.

5) The “chord problem” for Legendrian submanifolds in a contact manifold. In the case of a Legendrian knot k in \mathbb{R}^3 endowed with its standard contact form $\alpha = dz - y dx$, Y. Chekanov [13] gave a “combinatorial” Floer homology related to the chords of k , that is to the double points of the xy -projection of k . In that way he succeeded in discovering Legendrian knots with the same topological type and Thurston-Bennequin invariant, but which are not contact isotopic. In [20] the Chekanov complex is related to a Floer complex of J -holomorphic curves in the symplectization $(\mathbb{R}^3 \times \mathbb{R}, \omega = d(e^\tau \alpha))$ in the spirit of the *contact homology* theory from [19].

6) More about contact structures in dimension 3. P. Ozvath and Z. Szabo invented a *Heegaard-Floer homology* for 3-manifolds equipped with a Spin^c -structure [39]. When a contact structure ξ (and the associated Spin^c -structure)

is given on the 3-manifold Y , they define an integer $c(\xi)$ coming from calculations of some Heegaard-Floer groups. This invariant is non trivial but vanishes when ξ is *overtwisted*. So it is a deep new tool for studying *tight* contact structures.

Appendix : Topology of $\mathrm{Sp}(2n)$

A1 Local properties

We consider $\mathbb{R}^{2n} = \mathbb{C}^n$ endowed with its standard symplectic structure ω_0 and its complex structure J . The basis is ordered as a real basis of $\mathbb{C}^n = (\mathbb{R} + i\mathbb{R})^n$. The linear symplectic group $\mathrm{Sp}(2n)$ is the group of linear isomorphisms of $(\mathbb{R}^{2n}, \omega_0)$. It is identified with a group of invertible $2n \times 2n$ matrices. For $A \in \mathrm{Sp}(2n)$ close to I , one can write $A = e^{JS}$, where S is symmetric. Then $\mathrm{Sp}(2n)$ is a Lie group which has the same dimension as the space of quadratic forms on \mathbb{R}^{2n} , that is $n(2n+1)$. We recall notation from section 5:

$$\begin{aligned} \text{the Maslov cycle } \Sigma &:= \{A \in \mathrm{Sp}(2n) \mid \det(I - A) = 0\}, \\ \text{and } \mathrm{Sp}^*(2n) &:= \{A \in \mathrm{Sp}(2n) \mid \det(I - A) \neq 0\}. \end{aligned}$$

The Maslov cycle is stratified by the dimension of $V_1(A)$, where $V_1(A)$ denotes the eigenspace of A of the eigenvalue 1: $\Sigma = \Sigma_1 \sqcup \Sigma_2 \sqcup \cdots \sqcup \Sigma_{2n}$, where the union is made of disjoint strata

$$\Sigma_k = \mathrm{Sp}_k(2n) := \{A \in \mathrm{Sp}(2n) \mid \dim V_1(A) = k\}.$$

Since $A \mapsto \dim(V_1(A))$ is semi-continuous, the closure of Σ_k is $\overline{\Sigma}_k = \Sigma_k \cup \Sigma_{k+1} \cup \cdots \cup \Sigma_{2n}$. We recall that a linear isomorphism A of \mathbb{R}^{2n} is symplectic iff its graph Γ_A in $\mathbb{R}^{2n} \times \mathbb{R}^{2n}$ is Lagrangian with respect to $\omega_0 \oplus \omega_0$. The intersection of Γ_A with the diagonal Δ projects onto $V_1(A)$.

For a vector space E , $\mathrm{Sym}(E)$ denotes the space of self-adjoint operators $S : E \rightarrow E^*$ and $\mathrm{Sym}_k(E)$ denotes the subspace of those whose kernel has dimension k .

Proposition A1.1.

- 1) Σ_k is a smooth submanifold of codimension $\frac{k(k+1)}{2}$ in $\mathrm{Sp}(2n)$.
- 2) Let $A_0 \in \Sigma_k$ and $E_0 := \Gamma_{A_0} \cap \Delta$. There exists a germ of smooth embedding

$$T : (\mathrm{Sym}(E_0), 0) \rightarrow (\mathrm{Sp}(2n), A_0)$$

with the following properties:

- i) T is transversal to Σ_k at A_0 .
- ii) $T(s) \in \Sigma_j \iff s \in \mathrm{Sym}_j(E_0)$.

Proof. 1) In order to study $\mathrm{Sp}(2n)$ near A_0 we choose a Lagrangian complement Δ^* of the diagonal Δ transversal to Δ and to the graph Γ_{A_0} . It is indeed canonically isomorphic to the dual of Δ . Every Lagrangian subspace in $\mathbb{R}^{2n} \times \mathbb{R}^{2n}$ close to Γ_{A_0} is a Γ_A for some $A \in \mathrm{Sp}(2n)$.

Let A_0 be in Σ_k . A nearby A lies in Σ_k iff $\dim(\Gamma_A \cap \Delta) = k$. Let \mathcal{L} be the space of Lagrangian subspaces in $\Delta \oplus \Delta^*$ transversal to $0 \oplus \Delta^*$. Let us denote

$$\mathcal{L}_k := \{L \in \mathcal{L} \mid \dim(L \cap \Delta) = k\}.$$

Every $L \in \mathcal{L}$ has a *generating* quadratic form q . Precisely, L is the graph of a linear self-adjoint operator $B : \Delta \rightarrow \Delta^*$ and $q(x) = \langle B(x), x \rangle$. Let us split Δ into $\Delta = E_0 \oplus E_0'$ and dually $\Delta^* = E_0^* \oplus E_0'^*$. Then B reads as a matrix of operators

$$B = \begin{pmatrix} a & {}^t b \\ b & c \end{pmatrix}$$

where a and c are self-adjoint. For $L_0 := \Gamma_{A_0}$, the corresponding matrix is

$$B_0 = \begin{pmatrix} 0 & 0 \\ 0 & c_0 \end{pmatrix}$$

where c_0 is invertible. When B is close to B_0 , c is also invertible and

$$L \in \mathcal{L}_k \iff a + {}^t b c^{-1} b = 0 \quad (*).$$

The last equation says that the kernel of B has dimension k (and not less). Since this equation is regular near B_0 , we have proven 1). Moreover, in these coordinates, the equation of the tangent space $T_{L_0}(\mathcal{L}_k)$ reads $a = 0$.

2) As a consequence of the last remark the affine space of generating quadratic forms

$$\begin{pmatrix} a & 0 \\ 0 & c_0 \end{pmatrix} \quad (**)$$

is transversal to \mathcal{L}_k at Γ_{A_0} . Therefore when a is so small that L is a Γ_A , the corresponding family of A 's is transversal to Σ_k at A_0 . It is the transversal T we are looking for. Moreover, the dimension of the kernel of matrix $(**)$ is equal to the dimension of $\ker a$. Then 2ii) is clear. \square

From the above properties we immediately deduce the following corollary.

Corollary A1.2.

- 1) *The codimension of Σ_k in $\text{Sp}(2n)$ is at least 3 when $k > 1$.*
- 2) *$\overline{\Sigma}_2$ locally separates the Maslov cycle.*
- 3) *$\Sigma \setminus \overline{\Sigma}_3$ is locally connected.*

A2. Global connectivity properties

Lemma A2.1. $\Sigma_1 \cup \Sigma_2$ *is arcwise-connected.*

Proof. According to 3) from corollary A1.2, it is sufficient to prove that Σ is arcwise connected.

Let A_1 and A_2 be two symplectic isomorphisms with 1 as eigenvalue. Let us choose v_1 and v_2 two eigenvectors of A_1 and A_2 respectively. Because $\text{Sp}(2n)$ acts transitively on $\mathbb{R}^{2n} \setminus 0$ one can move A_2 in Σ by means of a 1-parameter family of conjugations in order to reduce ourselves to $v_2 = v_1 =: v$. The subgroup $\text{Sp}(2n)_v$ stabilizing v acts transitively on the set of Lagrangian subspaces

containing v . Then we can move A_1 and A_2 in $\mathrm{Sp}(2n)_v$ such that they share a common invariant Lagrangian subspace L . The subgroup $\mathrm{Sp}(2n)_L$ stabilizing L acts transitively on the set of Lagrangian subspaces transversal to L . Then we are reduced to the case where A_1 and A_2 keep invariant two transversal Lagrangian subspaces L and L^* , with $v \in L$. In that case there are a basis B of L , containing v , and its dual basis B^* of L^* such that, for $i = 1, 2$, A_i reads in (B, B^*) as

$$A_i = \begin{pmatrix} 1 & \times & & 0 \\ 0 & \Phi_i & & \\ & & 1 & 0 \\ 0 & & \times & {}^t\Phi_i^{-1} \end{pmatrix}$$

Now connecting A_1 and A_2 in Σ is an exercise which is left to the reader whatever the signs of $\det \Phi_i$ may be. \square

Proposition A2.2. $\Sigma_0 := \mathrm{Sp}(2n) \setminus \Sigma$ has two connected components.

Proof. One already knows that there are at least two components, one for each sign of $\det(I - A)$. Because $\mathrm{Sp}(2n)$ is connected, the closure of each component of Σ_0 meets Σ_1 . Hence it is sufficient to prove that, for every path $\gamma : [0, 1] \rightarrow \Sigma_1 \cup \Sigma_2$, $\gamma(t)$ adheres to the same component of $\mathrm{Sp}^+(2n)$ (resp. $\mathrm{Sp}^-(2n)$) as $\gamma(0)$ does. This is obvious as long as $\gamma(t) \in \Sigma_1$. The only problem could come when $\gamma(t_0) \in \Sigma_2$, because in this case $\gamma(t_0)$ seems to adhere to three components of Σ_0 according to proposition A1.1 (compare figure 5). Generically the 2-dimensional eigenspace $V_1(\gamma(t_0))$ is symplectic. In this case $\gamma(t_0)$ is conjugated to a normal form

$$\begin{pmatrix} 1 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & D \end{pmatrix}$$

Using the loop

$$t \in [0, 1] \mapsto \begin{pmatrix} \cos 2\pi t & -\sin 2\pi t & 0 \\ \sin 2\pi t & \cos 2\pi t & \\ 0 & 0 & D \end{pmatrix}$$

we see that, among the three local components of Σ_0 near $\gamma(t_0)$, two of them belong to the same component of Σ_0 . \square

A3. About $\pi_1(\mathrm{Sp}(2n) \setminus \Sigma)$ and co-orientation

First we recall that, when $A \in \Sigma_1$, the eigenvalue 1 has an even multiplicity because its characteristic polynomial is symmetric. Generically on Σ_1 the multiplicity is 2 and a normal form of A up to conjugation is

$$\begin{pmatrix} 1 & \lambda & 0 \\ 0 & 1 & \\ 0 & 0 & \Phi \end{pmatrix}.$$

Definition A3.1. We define

$$\text{sign}(A) := \text{sign}(\lambda \cdot \det(I - \Phi)).$$

Proposition A3.2. *The sign is well defined and locally constant on Σ_1 .*

Proof. Here is a simple proof due to N. Depauw and improved by J. Barge. I am pleased to thank both of them for this. Let $e_1(A)$ be the unique eigenvector of A with 1 as eigenvalue. Let us consider the endomorphism $M = I - A$. Since $A \in \Sigma_1$, the kernel of M is generated by $e_1(A)$ and the induced map on the quotient $M' : \mathbb{R}^{2n}/\mathbb{R}e_1(A) \rightarrow \mathbb{R}^{2n}$ is injective. Let H be the ω_0 -orthogonal of $e_1(A)$. As the ω_0 -product with $e_1(A)$ is preserved by A , we have $v - A(v) \in H$ for every $v \in \mathbb{R}^{2n}$. Hence the image of M' is a subspace of H . Due to the dimensions, M' is an isomorphism. The symplectic form allows us to identify the quotient $\mathbb{R}^{2n}/\mathbb{R}e_1(A)$ and H . Therefore M' can be seen as an isomorphism of H . Its determinant $d(A)$ is a well defined continuous function on Σ_1 . Using the above notation, we have $d(A) = \lambda \det(I - \Phi)$ when A is generic. \square

Notation A3.3.

- 1) Let us denote $\Sigma_1^\pm := \{A \in \Sigma_1 \mid \text{sign}(A) = \pm\}$ and $\Sigma^\pm := \overline{\Sigma_1^\pm} = \Sigma_1^\pm \cup \Sigma_2 \cup \dots \cup \Sigma_{2n}$.
- 2) Σ^+ is a hypersurface with singularities of codimension greater than 2 and it is co-oriented, the positive side being $\det(I - A) > 0$. Then Σ^+ is a homological cycle, having a well defined *intersection number* $[\Sigma^+].[\gamma]$ with every free homotopy class $[\gamma]$ of loops in $\text{Sp}(2n)$. When γ is in general position, the intersection number is the number of positive crossings minus the number of negative crossings.
- 3) For $k \in \mathbb{Z}$, let us define the loop γ_k by:

$$t \in [0, 1] \mapsto \gamma_k(t) = \begin{pmatrix} \cos 2\pi kt & -\sin 2\pi kt & 0 \\ \sin 2\pi kt & \cos 2\pi kt & 0 \\ 0 & 0 & D \end{pmatrix}.$$

Proposition A3.4. *Let γ be a loop in $\text{Sp}^\pm(2n)$. Then γ is homotopic to 0 in $\text{Sp}(2n)$.*

Proof. Since γ and Σ^+ are mutually disjoint we have $[\Sigma^+].[\gamma] = 0$. But we recall from the polar decomposition that γ is freely homotopic to the loop $t \mapsto M \cdot \gamma_k(t)$ for some $k \in \mathbb{Z}$, where M is any symplectic definite positive symmetric matrix independent of t , for instance:

$$M = \text{diag} \left(\begin{pmatrix} \xi & 0 \\ 0 & \xi^{-1} \end{pmatrix}, \dots, \begin{pmatrix} \xi & 0 \\ 0 & \xi^{-1} \end{pmatrix} \right), \quad \xi > 1.$$

If ξ is close to 1, an explicit calculation shows that $M \cdot \gamma_k$ crosses Σ^+ k times in the same direction. Thus $[\Sigma^+] \cdot [\gamma_k] = \pm k$, which forces the k associated with the given γ to be 0. (Compare figure 10). \square

Remark A3.5. Co-orientation of Σ

It is clear from the previous discussion that, if the smooth part of Σ is globally co-oriented by the sign of $\det(I - A)$, then Σ becomes a boundary in homology. For this reason we choose for Σ^- the opposite co-orientation: the *positive* side of Σ^- will be given by $\det(I - A) < 0$. With this co-orientation Σ^- and Σ^+ are homologous in $\text{Sp}(2n)$. As a consequence we have a clear geometric reason why the Maslov index of a loop is even.

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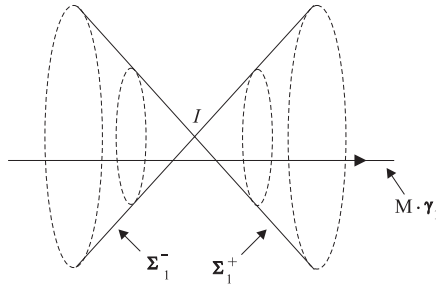


Figure 10

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François Laudenbach
Laboratoire Jean Leray, UMR 6629 du CNRS
Faculté des Sciences et Techniques Université de Nantes
F-44322 Nantes cedex 3, France
`laudenba@math.univ-nantes.fr`