Spectrum of the Second Variation

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Abstract

Second variation of a smooth optimal control problem at a regular extremal is a symmetric Fredholm operator. We study asymptotics of the spectrum of this operator and give an explicit expression for its determinant in terms of solutions of the Jacobi equation. In the case of the least action principle for the harmonic oscillator we obtain a classical Euler identity $\prod_{n=1}^{\infty} \left(1 - \frac{x^2}{(\pi n)^2}\right) = \frac{\sin x}{x}.$ General case may serve as a rich source of new nice identities.

1 Introduction

This research was initially motivated by the analysis of asymptotic expansions related to semigroups generated by hypo-elliptic operators. These expansions can be often interpreted as infinite-dimensional versions of the standard asymptotics for expressions of the form $\int\limits_{\mathbb{R}^N} a(u)e^{\frac{\varphi(u)}{-t}}du \text{ or } \int\limits_{\mathbb{R}^N} a(u)e^{\frac{\varphi(u)}{it}}du$ as $t\to 0$. It is well-known that the asymptotics is localized in the critical points x_i of φ , and the quantities $(\det D^2_{x_i})^{-\frac{1}{2}}$ play a crucial role in the expansions. In the infinite-dimensional version we deal with a "path integral", and φ is an "action functional" whose critical points are extremals of the appropriate optimal control problem. So we try to analyse and compute $\det D^2_{x_i}$ in this case. What we obtain, can be considered as a modification and generalization of the classical "Hill's determinant" in Rational Mechanics (see [2]).

Before the calculation of the determinant, we study asymptotics of the spectrum of the second variation $D_u^2 \varphi$ that is a symmetric Fredholm operator of the form I + K, where K is a compact Hilbert–Schmidt operator. The

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point is that K is usually NOT a trace class operator so that the trace of K and the determinant of I + K are not well-defined in the standard sense.

Anyway, a specific structure of the spectral asymptotics allows to define and compute everything. This asymptotic is described in Theorem 1. Actually, there is an important symmetry in the asymptotics that leads to a cancellation of slow convergent to zero positive and negative terms when we compute the sum or the product.

There is a general reason for this symmetry. It concerns an evolutionary nature of the optimal control problems: the space of available control functions grows with time. This evolutionary structure results in the following property (I would call it "causality") of the operator K: there exists a Volterra operator V on $L^2([0,1];\mathbb{R}^m)$ and a finite codimension subspace of $L^2([0,1];\mathbb{R}^m)$ such that $\langle Ku,u\rangle=\langle Vu,u\rangle$ for any u from the subspace. In other words, our symmetric operator imitates a triangle one on a subspace of finite codimension. We expect that all causal operators have spectral asymptotics similar to the asymptotics from Theorem 1, although we do not study general causal operators in this paper.

In Theorem 2 we give an explicit expression for $\det(I+K)$ in terms of solutions of the Jacobi equation that generalize classical Jacobi fields. We also give an explicit integral expression for $\operatorname{tr} K$; moreover, integral formulas for all elementary symmetric functions of the eigenvalues of K can be recovered from the main determinantal formula.

I hope that this result would serve as an effective summation method, a way to get explicit expressions for infinite products and sums of interesting series. A simple example: for the 1-dimensional linear control system $\dot{x} = ax + u \text{ with the quadratic cost } \varphi(u) = \int_0^1 u^2(t) - (a^2 + b^2) x^2(t) \, dt \text{ our determinantal identity reads: } \prod_{n=1}^\infty \left(1 - \frac{a^2 + b^2}{a^2 + (\pi n)^2}\right) = \frac{a \sin b}{b \sin a}; \text{ the case } a = 0$ corresponds to the famous Euler identity $\prod_{n=1}^\infty \left(1 - \frac{b^2}{(\pi n)^2}\right) = \frac{\sin b}{b}.$

2 Preliminaries

We consider a smooth control system of the form:

$$\dot{q} = f_u^t(q), \quad q \in M, \ u \in U, \tag{1}$$

with a fixed initial point $q_0 \in M$. Here M and U are smooth manifolds (without border) and the vector $f_u^t(q) \in T_qM$ smoothly depends on $(q, u) \in M \times U$ and is measurable bounded with respect to $t \in [0, 1]$.

Let $u(\cdot):[0,1]\to U$ be a measurable map such that u([0,1]) is contained in a compact subset of U. The Cauchy problem

$$\dot{q} = f_{u(t)}^t(q), \quad q(0) = q_0,$$

has a unique Lipschitzian solution $t \mapsto q(t; u(\cdot))$ defined on an interval in \mathbb{R} . We say that $u(\cdot)$ is an admissible control and $q(\cdot; u(\cdot))$ is a corresponding admissible trajectory if the domain of this solution contains [0,1]. We denote the set of all admissible controls by \mathcal{U} ; then \mathcal{U} is an open subset of $L^{\infty}([0,1]; \mathcal{U})$. Hence \mathcal{U} is a smooth Banach manifold modelled on the space $L^{\infty}([0,1]; \mathbb{R}^{\dim \mathcal{U}})$.

Given $t \in [0,1]$ we define the "evaluation map" $F_t : \mathcal{U} \to M$ by the formula $F_t(u(\cdot)) = q(t;u(\cdot))$; then F_t is a smooth map from the Banach manifold \mathcal{U} to M.

Let $\ell: M \times U \to \mathbb{R}$ be a smooth "Lagrangian". We consider functionals $\varphi_t: \mathcal{U} \to \mathbb{R}, \ 0 \le t \le 1$, defined by the formula:

$$\varphi_t(u(\cdot)) = \int_0^t \ell(q(\tau; u(\cdot)), u(\tau)) d\tau.$$
 (2)

Definition 1. We say that $u \in \mathcal{U}$ is a normal extremal control if there exists $\lambda_1 \in T_{F_1(u)}^*M$ such that $\lambda_1 D_u F_1 = d_u \varphi_1$; here $\lambda_1 D_u F_1$ is the composition of $D_u F_1 : T_u \mathcal{U} \to T_{F_1(u)} M$ and $\lambda_1 : T_{F_1(u)} M \to \mathbb{R}$. We say that a normal extremal control is strictly normal if it is a regular point of F_1 .

Strictly normal extremal controls are just critical points of $\varphi|_{F_1^{-1}(q)}$, $q \in M$. Let $t \in [0,1]$; it is easy to see that the restriction $u|_{[0,t]}$ of a normal extremal control u is also a normal extremal control in the following sense:

$$\exists \lambda_t \in T^*_{F_t(u)} M \quad \text{such that} \quad \lambda_t D_u F_t = d_u \varphi_t. \tag{3}$$

The "differentiation of the identity (3) with respect to t" leads to the following Hamiltonian characterization of normal extremal controls.

A family of Hamiltonians $h_u^t: T^*M \to \mathbb{R}, \ u \in U$, is defined by the formula:

$$h_u^t(\lambda) = \langle \lambda, f_u^t(q) \rangle - \ell(q,u), \quad g \in M, \ \lambda \in T_q^*M.$$

Let σ be the canonical symplectic form on T^*M and $\pi: T^*M \to M$ be the standard projection, $\pi(T_q^*M) = q$. Recall that $\sigma = ds$ where s is a Liouville (or tautological) 1-form, $\langle s_{\lambda}, \eta \rangle = \langle \lambda, \pi_* \eta \rangle$, $\forall \lambda \in T^*M$, $\eta \in T_{\lambda}(T^*M)$. Given a smooth function $h: T^*M \to \mathbb{R}$, the Hamiltonian vector field \vec{h} on T^*M is defined by the identity: $dh = \sigma(\cdot, \vec{h})$.

Proposition 1. Let $\tilde{u} \in \mathcal{U}$ and $\tilde{q}(t) = q(t; \tilde{u}), \ 0 \le t \le 1$; then \tilde{u} is a normal extremal control if and only if there exists $\tilde{\lambda}_t \in T^*_{\tilde{q}(t)}M$ such that

$$\dot{\tilde{\lambda}}_t = \vec{h}_{\tilde{u}(t)}^t, \quad \frac{\partial h_u^t(\tilde{\lambda}_t)}{\partial u}\Big|_{u=\tilde{u}(t)} = 0, \quad 0 \le t \le 1.$$
(4)

This statement is standard, it is actually a weak version of the Pontryagin maximum principle.

A Lipschitzian curve $\tilde{\lambda}_t$, $0 \leq t \leq 1$, in T^*M that satisfies relations (4) is called a normal extremal associated to the normal extremal control \tilde{u} . The time-varying Hamiltonian system $\dot{\lambda} = \vec{h}_{\tilde{u}(t)}^t(\lambda)$ defines a flow

$$\tilde{\Phi}^t: T^*M \to T^*M, \quad \tilde{\Phi}^t: \lambda(0) \mapsto \lambda(t),$$

where $\dot{\lambda}(\tau) = \vec{h}_{\tilde{u}(\tau)}^t(\lambda(\tau)), \ 0 \le \tau \le t.$ Obviously, $\tilde{\Phi}^t(\tilde{\lambda}_0) = \tilde{\lambda}_t$. Moreover, $\pi_* \vec{h}_{\tilde{u}(t)}^t(\lambda) = f_{\tilde{u}(t)}^t(q), \ \forall \, q \in M, \ \lambda \in T_q^*M$.

Let $\tilde{P}^t: M \to M$ be the flow generated by the time-varying system $\dot{q} = f^t_{\tilde{u}(t)}(q)$, i.e. $\tilde{P}^t: q(0) \mapsto q(t)$, where $\dot{q}(\tau) = f^\tau_{\tilde{u}(\tau)}(q(\tau))$, $0 \le \tau \le t$. It follows that $\tilde{\Phi}^t$ are fiberwise transformations and $\tilde{\Phi}^t(T^*_qM) = T^*_{\tilde{P}^t(q)}M$. Moreover, the restriction of $\tilde{\Phi}^t$ to T^*_qM is an affine map of the vector space T^*_qM on the vector space $T^*_{\tilde{P}^t(q)}M$.

In what follows, we assume that \tilde{u} is a normal extremal control and λ_t is a normal extremal, they are fixed until the end of the paper. Then $\tilde{u}(t)$ is a critical point of the function $u\mapsto h^t_u(\tilde{\lambda}_t),\ u\in U$ for any $t\in [01]$. Hence the Hessian $\frac{\partial^2 h^t_{\tilde{u}(t)}}{\partial u^2}(\tilde{\lambda}_t)$ is a well-defined quadratic form on $T_{\tilde{u}(t)}U$ or, in other words, a self-adjoint linear map from $T_{\tilde{u}(t)}U$ to $T^*_{\tilde{u}(t)}U$.

Now we consider Hamiltonian functions

$$g_u^t = (h_u^t - h_{\tilde{u}(t)}^t) \circ \tilde{\Phi}^t, \quad u \in U, \ t \in [0, 1].$$

A time-varying Hamiltonian vector field \vec{g}_u^t generates the flow $(\tilde{\Phi}^t)^{-1} \circ \Phi_u^t$, where Φ_u^t is the flow generated by the field \vec{h}_u^t . We have:

$$g_{\tilde{u}(t)}^t \equiv 0, \quad \frac{\partial g_{\tilde{u}(t)}^t}{\partial u}(\tilde{\lambda}_0) = 0, \quad \frac{\partial^2 g_{\tilde{u}(t)}^t}{\partial^2 u}(\tilde{\lambda}_0) = \frac{\partial^2 h_{\tilde{u}(t)}^t}{\partial^2 u}(\tilde{\lambda}_t).$$

We introduce a simplified notation $H_t \doteq \frac{\partial^2 h_{\tilde{u}(t)}^t}{\partial^2 u}(\tilde{\lambda}_t)$; recall that $H_t : T_{\tilde{u}(t)}U \to T_{\tilde{u}(t)}^*U$ is a self-adjoint linear map.

More notations. Consider the map $u \mapsto \vec{g}_u^t(\tilde{\lambda}_0)$ from U to $T_{\tilde{\lambda}_0}(T^*M)$ We denote by Z_t the differential of this map at $\tilde{u}(t)$, $Z_t \doteq \frac{\partial \vec{g}_{\tilde{u}(t)}^t}{\partial u}(\tilde{\lambda}_0)$; then Z_t is a linear map from $T_{\tilde{u}(t)}U$ to $T_{\tilde{\lambda}_0}(T^*M)$. We also set $X_t = \pi_*Z_t$; then X_t is a linear map from $T_{\tilde{u}(t)}U$ to $T_{\tilde{q}(t)}M$. Finally, we denote by $J: T_{\tilde{\lambda}_0}(T^*M) \to T_{\tilde{\lambda}_0}^*(T^*M)$ the anti-symmetric linear map defined by the identity $\sigma_{\tilde{\lambda}_0}(\cdot,\cdot) = \langle J\cdot,\cdot\rangle$.

Note that $T_{u(\cdot)}\mathcal{U}$ is the space of measurable bounded mappings $t \mapsto v(t) \in T_{u(t)}U$, $0 \le t \le 1$. A simple rearrangement of the standard formula for the first variation in the optimal control theory (see for instance the textbook [1]) gives the identity:

$$(D_{\tilde{u}}F_t)v = \tilde{P}_*^t \int_0^t X_\tau v(\tau) d\tau.$$

Assume that \tilde{u} is a regular point of F_t (i. e. $D_{\tilde{u}}F_t$ is surjective); then \tilde{u} is a critical point of $\varphi_t|_{F_t^{-1}(\tilde{q}(t))}$. We have:

$$T_{\tilde{u}}(F_t^{-1}(\tilde{q}(t)) = \ker D_{\tilde{u}}F_t = \left\{ v \in T_{\tilde{u}}\mathcal{U} : \int_0^t X_{\tau}v(\tau) d\tau = 0 \right\}.$$

The Hessian of $\varphi_1|_{F_1^{-1}(\tilde{q}(1))}$ at \tilde{u} is a quadratic form $D^2\varphi_1$: ker $D_{\tilde{u}}F_1 \to \mathbb{R}$. This is the "second variation" of the optimal control problem at \tilde{u} , the main object of this paper. It has the following expression (see [1]):

$$D_{\tilde{u}}^{2}\varphi_{1}(v) = -\int_{0}^{1} \langle H_{t}v(t), v(t)\rangle dt - \int_{0}^{1} \left\langle J \int_{0}^{t} Z_{\tau}v(\tau) d\tau, Z_{t}v(t) \right\rangle dt, \quad (5)$$

where $v \in T_{\tilde{u}}\mathcal{U}$ and $\int_0^1 X_t v(t) dt = 0$.

Definition 2. The extremal $\tilde{\lambda}_t$, $0 \le t \le 1$, is regular if H_t is invertible for any $t \in [0, 1]$ with a uniformly bounded inverse.

In what follows, we assume that the reference extremal is regular. Moreover, we assume that $H_t < 0^1$ for any $t \in [0,1]$. This last assumption is motivated by the classical Legendre condition: for a regular extremal, integral quadratic form (5) has a finite negative inertia index if and only if $H_t < 0$, for any $t \in [0,1]$.

Finally we introduce a "Gramm matrix", a self-adjoint linear map Γ_t : $T_{q_0}^*M \to T_{q_0}M$ defined by the formula: $\Gamma_t = -\int_0^t X_\tau H_t^{-1} X_\tau^* d\tau$. We see that \tilde{u} is a regular point of F_t if and only if Γ_t is invertible.

¹We say that a self-adjoint map is positive (negative) if the corresponding quadratic form is positive (negative).

Remark 1. There is an apparently more general and natural way to define a control system. Let me briefly describe it confining within the time-invariant case. Indeed, instead of the product $M \times U$ we may consider a locally trivial bundle over M with a typical fiber U. Then f is a smooth fiberwise map from this bundle to TM that sends the fiber U_q into T_qM , $q \in M$. In this setting, the control and the correspondent trajectory is somehow a unique object: $u(t) \in U_{q(t)}$ and $\dot{q}(t) = f(u(t))$.

The situation is reduced to what we did before this remark if the bundle is trivial. In the general case, we simply trivialize the bundle in a neighborhood of the reference trajectory. To be more precise, we first take $\mathbb{R} \times M$, where \mathbb{R} is the time axis and trivialize the pullback bundle over $\mathbb{R} \times M$ in a neighborhood of the graph of the reference trajectory (the graph does not have self-intersections while the original trajectory might have). Moreover, all statements of this paper use only H_t , Z_t , and X_t , $0 \le t \le 1$, and these quantities depend only on the trivialization of the vector bundle $T_{\tilde{u}(t)}U_{q(t)}$ along (the graph of) the curve u(t), $0 \le t \le 1$. A trivialization if the vector bundle along such a curve is naturally achieved by a parallel transport in virtue of a linear connection on the vector bundle with the fibers T_uU_q , $q \in M$, $u \in U_q$ and the base $\bigcup_{q \in M} U_q$.

So, a linear connection on this vector bundle is actually all we need to write all the formulas. I leave to an interested reader to write them. I have decided not to follow this way in order to avoid an unnecessary language complication and to make the paper affordable for a larger audience.

3 Main results

According to our construction, the space $T_{\tilde{u}}\mathcal{U}$ consists of the L^{∞} -maps $t \mapsto v(t) \in T_{\tilde{u}(t)}U$, $0 \leq t \leq 1$. At the same time, linear map $D_{\tilde{u}}F_1: T_{\tilde{u}}\mathcal{U} \to T_{\tilde{q}(1)}M$ and quadratic form $D_{\tilde{u}}^2\varphi_1$ are continuous in a weaker topology L^2 . Let \mathcal{V} be the closure of ker $D_{\tilde{u}}F_1$ in the topology L^2 . Then \mathcal{V} is a Hilbert space equipped with a Hilbert structure

$$< v_1 | v_2 > \doteq \int_0^1 \langle -H_t v_1(t), v_2(t) \rangle dt.$$

Formula (5) implies that

$$D_{\tilde{u}}^{2}\varphi_{1}(v) = \langle (I+K)v|v\rangle, \quad v \in \mathcal{V}, \tag{6}$$

where K is a compact symmetric operator on \mathcal{V} . In particular, the spectrum

of K is real, the only limiting point of the spectrum is 0, and any nonzero eigenvalue has a finite multiplicity.

What can we say about the "trace" of K and the "determinant" of I+K? Let us consider a simple example, the least action principle for a charged particle in the plane in a constant magnetic field.

Example 1. Let $M=U=\mathbb{R}^2,\ q=(q^1,q^2),\ u=(u^1,u^2),\ f_u(q)=u,\ \ell(q,u)=\frac{1}{2}|u|^2+r(q^1u^2-q^2u^1),\ q_0=0,\ \tilde{u}(t)=0.$ Then

$$T^*M = \mathbb{R}^2 \times \mathbb{R}^2 = \{(p,q) : p \in \mathbb{R}^2, q \in \mathbb{R}^2\}, J(p,q) = (-q,p),$$

$$h_u(p,q) = \langle p, u \rangle - \frac{1}{2}|u|^2 - r(q^1u^2 - q^2u^1),$$

$$h_{\tilde{u}(t)} = 0, \quad \tilde{\lambda}_t = (\tilde{p}(t), \tilde{q}(t)) = (0, 0), \quad g_u^t = h_u.$$

We have $H_t v = -v$, $Z_t v = (rv^2, -rv^1; v^1, v^2)$, $X_t v = v$, $t \in [0, 1]$, $v \in \mathbb{R}^2$. Then

$$\mathcal{V} = \{ v \in L^2([0,1]; \mathbb{R}^2) : \int_0^1 v(t) \, dt = 0 \}.$$

It is convenient to identify \mathbb{R}^2 with \mathbb{C} as follows: $(v^1, v^2) = v^1 + iv^2$. An immediate calculation gives the following expression for the operator K:

$$Kv(t) = \int_0^t 2riv(\tau) d\tau - \int_0^1 \int_0^t 2riv(\tau) d\tau dt.$$

The eigenfunctions of this operator have a form $t \mapsto ce^{2\pi nit}$, $0 \le t \le 1$, $c \in \mathbb{C}$, $n = \pm 1, \pm 2, \ldots$, where the eigenfunction $ce^{2\pi nit}$ corresponds to the eigenvalue $\frac{r}{\pi n}$.

We see that even in this model example the eigenvalues of K do not form an absolutely convergent series and K is not a trace class operator. On the other hand, the next theorem implies that a "principal value" of such a series does exist at least if the data are piece-wise real analytic.

Consider the operator $H_t^{-1}Z_t^*JZ_t: T_{\tilde{u}(t)}U \to T_{\tilde{u}(t)}U$. This operator is associated to an anti-symmetric bilinear form and has only purely imaginary eigenvalues. We denote by $\bar{\zeta}_t$ the sum of positive eigenvalues of the operator $iH_t^{-1}Z_t^*JZ_t$ counted according to multiplicity (here i is the imaginary unit).

Remark 2. Analysing notations of Section 2 we see that $\bar{\zeta}_t$ is the sum of positive roots (counted according to multiplicity) of the equation

$$\det\left(\left\{\frac{\partial h_{\tilde{u}(t)}^t}{\partial u}, \frac{\partial h_{\tilde{u}(t)}^t}{\partial u}\right\}(\tilde{\lambda}_t) + si\frac{\partial^2 h_{\tilde{u}(t)}^t}{\partial u^2}(\tilde{\lambda}_t)\right) = 0$$

with unknown s. Here $\{\cdot,\cdot\}$ is the Poisson bracket so that

$$(v,v') \mapsto \left\{ \frac{\partial h_{\tilde{u}(t)}^t}{\partial u} v, \frac{\partial h_{\tilde{u}(t)}^t}{\partial u} v' \right\} (\tilde{\lambda}_t), \quad v,v' \in T_{\tilde{u}(t)} U,$$

is an anti-symmetric bilinear form and $\frac{\partial^2 h_{\tilde{u}(t)}^t}{\partial u^2}(\tilde{\lambda}_t)$, is a symmetric bilinear form on $T_{\tilde{u}(t)}U$.

Let $\operatorname{Sp}(K) \subset \mathbb{R}$ be the spectrum of the operator K, $\operatorname{Sp}(K) \setminus \{0\} = \operatorname{Sp}_{+}(K) \cup \operatorname{Sp}_{-}(K)$, where $\operatorname{Sp}_{\pm}(K) \subset \mathbb{R}_{\pm}$. Given $\alpha \in \operatorname{Sp}(K) \setminus 0$, we denote by m_{α} the multiplicity of the eigenvalue α . Moreover, if $\operatorname{Sp}_{\pm}(K)$ is an infinite set, then we introduce a natural ordering of $\operatorname{Sp}_{\pm}(K)$ that is a monotone decreasing sequence α_n , $n \in \mathbb{Z}_{\pm}$, with the following properties:

$$\bigcup_{n \in \mathbb{Z}_{\pm}} \{\alpha_n\} = \operatorname{Sp}_{\pm}(K), \quad \#\{n \in \mathbb{Z}_{\pm} : \alpha_n = \alpha\} = m_{\alpha}. \tag{7}$$

Theorem 1. If $\bar{\zeta}_t \equiv 0$, then $\alpha_n = O(|n|^{-2})$ as $n \to \pm \infty$. If $\bar{\zeta}_t$ is not identical zero and H_t and Z_t are piecewise real analytic with respect to t then $\mathrm{Sp}_+(K)$ and $\mathrm{Sp}_-(K)$ are both infinite and

$$\alpha_n = \frac{1}{\pi n} \int_0^1 \bar{\zeta}_t \, dt + O\left(|n|^{-5/3}\right) \quad \text{as} \quad n \to \pm \infty.$$
 (8)

Remark 3. It is reasonable to expect that the statement of the theorem is valid without the piecewise-analyticity assumption. Moreover, the order $n^{-5/3}$ of the remainder term is certainly not optimal and perhaps can be substituted by n^{-2} . Anyway, the stated result is quite sufficient for our purposes while the proof of a stronger one would require a more sophisticated technique.

A cancellation of slow convergent to zero terms of the opposite sign in the expansion (8) gives the following:

Corollary 1. The depending on $\varepsilon > 0$ families of real numbers

$$\sum_{\substack{\alpha \in \operatorname{Sp}(K) \\ |\alpha| \ge \varepsilon}} m_{\alpha} \alpha, \quad \prod_{\substack{\alpha \in \operatorname{Sp}(K) \\ |\alpha| \ge \varepsilon}} (1+\alpha)^{m_{\alpha}}$$

have finite limits as $\varepsilon \to 0$.

We use natural notations for these limits:

$$\operatorname{tr} K = \lim_{\varepsilon \to 0} \sum_{\substack{\alpha \in \operatorname{Sp}(K) \\ |\alpha| \ge \varepsilon}} m_{\alpha} \alpha, \quad \det(I + K) = \lim_{\varepsilon \to 0} \prod_{\substack{\alpha \in \operatorname{Sp}(K) \\ |\alpha| \ge \varepsilon}} (1 + \alpha)^{m_{\alpha}}.$$

We are going to compute these trace and determinant in terms of H_t , Z_t , X_t and solutions of the following Jacobi system:

$$\dot{\eta} = -Z_{\tau} H_{\tau}^{-1} Z_{\tau}^* J \eta, \quad \eta(t) \in T_{\tilde{\lambda}_0}(T^* M), \ 0 \le \tau \le 1.$$
 (9)

This is a linear time-varying Hamiltonian system in $T_{\tilde{\lambda}_0}(T^*M)$ associated

to a nonnegative quadratic Hamiltonian $H_t(\eta) = -\frac{1}{2} \langle Z_t^* J \eta, H_t^{-1} Z_t^* J \eta \rangle$. In what follows, we identify the space $T_{q_0}^* M$ with its tangent $T_{\tilde{\lambda}_0}(T_{q_0}^* M) \subset T_{\tilde{\lambda}_0}(T^*M)$ and introduce linear maps $Q_t: T_{q_0}^* M \to T_{q_0} M, \ 0 \le t \le 1$, by the formula: $Q_t(\eta(0)) = \pi_* \eta(t)$, where $\eta(\tau)$, $0 \le \tau \le t$, is a solution of equation (9) and $\eta(0) \in T_{q_0}^* M$.

Remark 4. According to our assumptions, $\tilde{u}(\tau)$ is a strict local maximum of the function $u \mapsto h_u^{\tau}(\lambda_{\tau})$. Assume that this is a global maximum and moreover $h^{\tau}(\lambda) = \max_{u \in U} h_u^{\tau}(\lambda)$ is smooth with respect to $\lambda \in T^*M$. Then $\tilde{\lambda}_{\tau}, \ 0 \leq \tau \leq 1$, is a solution of the Hamiltonian system $\dot{\lambda} = \vec{h}^{\tau}(\lambda)$. We define the exponential map $\mathcal{E}_q^t: T_q^*M \to M$ by the formula $\mathcal{E}_q^t(\lambda_0) = \pi(\lambda_t)$, where $\dot{\lambda}_{\tau} = \vec{h}^{\tau}(\lambda_{\tau}), \ 0 \le \tau \le t.$ It is not hard to show that $Q_t = (\tilde{P}_*^t)^{-1} D_{\tilde{\lambda}_0} \mathcal{E}_{q_0}^t$.

Theorem 2. Under conditions of Theorem 1, the following identities are valid:

$$\det(I+K) = \det(Q_1 \Gamma_1^{-1}),$$

$$\operatorname{tr} K = \operatorname{tr} \left(\int_0^1 \int_0^t X_t H_t^{-1} Z_t^* J Z_\tau H_\tau^{-1} X_\tau^* d\tau dt \Gamma_1^{-1} \right).$$

Let us apply this theorem to Example 1. In our coordinates, $\Gamma_1 = I$. We have to find matrix Q_1 . It is convenient to use complex notations: $\eta = (p;q) \in \mathbb{C} \times \mathbb{C}, \ p = ip_1 + p_2, \ q = iq_1 + q_2$. System (8) has the form: $\dot{p} = irp - \nu^2 q$, $\dot{q} = irq + p$. We have to find q(t) under conditions p(0) = 1, q(0) = 0; then Q_1 is just complex number q(1) treated as a 2×2 real matrix, det $Q_1 = |q(1)|^2$. A simple calculation gives: $q(t) = \frac{\sin \nu t}{\nu} e^{i\nu t}$. Keeping in mind that all eigenvalues of the operator K have multiplicity 2, we obtain the square of a classical Euler identity: $\prod_{n=1}^{+\infty} \left(1 - \left(\frac{r}{\pi n}\right)^2\right)^2 =$ $\left(\frac{\sin r}{r}\right)^2$.

Now we consider one more very simple example, a harmonic oscillator.

Example 2. $M=U=\mathbb{R},\ f_u(q)=u,\ \ell(q,u)=\frac{1}{2}(u^2-rq^2),\ q_0=0,\ \tilde{u}(t)=0.$ Then $h_u(p,q)=pu-\frac{1}{2}(u^2-rq^2),\ h_{\tilde{u}(t)}=\frac{r}{2}q^2,$

$$\tilde{\Phi}_t(p,q) = \begin{pmatrix} p - trq \\ q \end{pmatrix}, \ g_u^t = (p - trq)u - \frac{1}{2}u^2, \quad H_t = -1, \ Z_t = \begin{pmatrix} tr \\ 1 \end{pmatrix}, \ X_t = 1.$$

Operator K has a form:

$$Kv(t) = r \int_0^t (t - \tau)v(\tau) d\tau - r \int_0^1 \int_0^t (t - \tau)v(\tau) d\tau dt.$$

The eigenfunctions of this operator have a form $t\mapsto c\cos(\pi nt),\ c\in\mathbb{R},\ n=1,2,\ldots$, where the eigenfunction $c\cos(\pi nt)$ corresponds to the eigenvalue $-\frac{r}{(\pi n)^2}$. Moreover, $Q_1=\frac{\sin\sqrt{r}}{\sqrt{r}}$ if r>0 and $Q_1=\frac{\sinh\sqrt{|r|}}{\sqrt{|r|}}$ if r<0. The determinant formula from Theorem 2 coincides with the Euler identity: $\prod_{n=1}^{\infty}\left(1-\frac{r}{(\pi n)^2}\right)=\frac{\sin\sqrt{r}}{\sqrt{r}} \text{ or its hyperbolic version.}$ The trace formula gives another famous Euler observation: $\sum_{n=1}^{\infty}\frac{r}{(\pi n)^2}=\frac{r}{6}.$

Harmonic oscillator is a special case of a more general example where the spectrum is still explicitly computed.

Example 3. $M = U = \mathbb{R}^m$, $f_u(q) = Aq + u$, $\ell(q, u) = \frac{1}{2}(|u|^2 - \langle q, Rq \rangle)$, $q_0 = 0$, $\tilde{u}(t) = 0$, where A is a $m \times m$ -matrix and R is a symmetric $m \times m$ -matrix. Then $h_u(p,q) = \langle p, Aq + u \rangle - \frac{1}{2}(|u|^2 - \langle q, Rq \rangle)$, $h_{\tilde{u}(t)} = \langle p, Aq \rangle + \frac{1}{2}\langle q, Rq \rangle$,

$$\tilde{\Phi}_{t}(p,q) = \begin{pmatrix} e^{-tA^{*}}p - \int_{0}^{t} e^{(\tau-t)A^{*}}Re^{\tau A} d\tau \\ e^{tA}q \end{pmatrix},
g_{u}^{t} = \langle e^{-tA^{*}}p - \int_{0}^{t} e^{(\tau-t)A^{*}}Re^{\tau A} d\tau, u \rangle - \frac{1}{2}|u|^{2},
H_{t} = -I, \quad Z_{t} = \begin{pmatrix} \int_{0}^{t} e^{\tau A^{*}}Re^{\tau A} d\tau \\ e^{-tA} \end{pmatrix}, \quad X_{t} = e^{-tA}.$$

Then $\mathcal{V} = \{ v \in L^2([0,1]; \mathbb{R}^2) : \int_0^1 e^{-tA} v(t) dt = 0 \},$

$$< Kv|v> = \int_0^1 \left\langle \iint_{0 \le \tau_1 \le \tau_2 \le t} e^{(\tau_2 - t)A^*} Re^{(\tau_2 - \tau_1)A} v(\tau_1) d\tau_2 d\tau_1, v(t) \right\rangle dt.$$

A vector function $v(\cdot) \in \mathcal{V}$ is an eigenfunction of the operator K with an eigenvalue $\alpha \in \mathbb{R}$ if and only if

$$\alpha v(t) = \iint_{0 \le \tau_1 \le \tau_2 \le t} e^{(\tau_2 - t)A^*} Re^{(\tau_2 - \tau_1)A} v(\tau_1) d\tau_2 d\tau_1 + e^{-tA^*} c, \quad t \in [0, 1], (10)$$

for some constant vector $c \in \mathbb{R}^m$. We denote: $y(t) = \int_0^t e^{(t-\tau)A} v(\tau) d\tau$; then we differentiate twice equation (10) and obtain that this equation is equivalent to the boundary values problem y(0) = y(1) = 0 for the ordinary differential equation

$$\alpha \ddot{y} = \alpha (A - A^*)\dot{y} + (\alpha A^*A + R)y.$$

From now on we assume that the matrix A is symmetric, $A^* = A$; then nonzero eigenvalues and corresponding eigenfunctions are nontrivial solution of the boundary value problem:

$$\ddot{y} = (A^2 + \frac{1}{\alpha}R)y, \quad y(0) = y(1) = 0.$$

We obtain that α is a nonzero eigenvalue of K if and only if there exists a positive integer n such that $-(\pi n)^2$ is an eigenvalue of the matrix $A^2 + \frac{1}{\alpha}R$. Moreover, the multiplicities of the eigenvalue α of K and of the eigenvalue $-(\pi n)^2$ of $A^2 + \frac{1}{\alpha}R$ are equal. In other words, $\det K = \prod_{n=1}^{\infty} \prod_i (1 - s_{ni})$ where s_{ni} are real roots (counted with their multiplicities) of the following polynomial equation with unknown s:

$$\det(s(A^{2} + (\pi n)^{2}I) - R) = 0.$$

Actually, all roots of this polynomial are real since $A^2 + (\pi n)^2 I$ is a signdefinite symmetric matrix. Let $\phi(s) = \det(s(A^2 + (\pi n)^2 I) - R)$, $\phi(s) = \sum_{i=1}^{m} c_{ni} s^i$; then $c_{nm} = \det(A^2 + (\pi n)^2 I)$ and $\prod_{n=1}^{\infty} \prod_{i=1}^{\infty} (1 - s_{ni}) = \frac{\phi(1)}{c_{nm}}$. We

$$\prod_{i=1}^{m} (1 - s_{ni}) = \frac{1}{c_{nm}} \det((A^2 + (\pi n)^2 I) - R) = \det(I - R(A^2 + (\pi n)^2 I)^{-1}).$$

Hence det
$$K = \prod_{n=1}^{\infty} \det(I - R(A^2 + (\pi n)^2 I)^{-1}).$$

Similarly,
$$\text{tr}K = -\sum_{n=1}^{\infty} \text{tr}(R(A^2 + (\pi n)^2 I)^{-1}).$$

Similarly, $\operatorname{tr} K = -\sum_{n=1}^{\infty} \operatorname{tr}(R(A^2 + (\pi n)^2 I)^{-1}).$ Now we have to compute the right-hand sides of the identities of Theorem 2. First of all, $\Gamma_1 = \int_0^1 e^{-2tA} dt$. Then we consider Jacobi system (9). We have $T^*M = \mathbb{R}^m \times \mathbb{R}^m = \{(p,q) : p,q \in \mathbb{R}^m\}$; then $Q_t p_0 = q_t$ where $t \mapsto (p_t, q_t)$ is the solution of the Jacobi system with the initial value $q_0 = 0$. We set $y(t) = e^{tA}q_t$, differentiate in virtue of the Jacobi system and obtain the equation: $\ddot{y} = (A^2 - R)y$. Moreover, $\dot{y}(0) = p_0$. Hence

 $Q_t = (R - A^2)^{-1/2} \sin((A^2 - R)^{1/2}t)$. This formula is valid also for a sign-indefinite matrix $R - A^2$ if we properly interpret the square root or simply make the computations in coordinates where the matrix $A^2 - R$ is diagonal. Finally, putting together all the formulas, we obtain the following generalization of the Euler identities:

Proposition 2. Let A, R be symmetric matrices. Then:

$$\prod_{n=1}^{\infty} \det \left(I - R(A^2 + (\pi n)^2 I)^{-1} \right) = \frac{2 \det \left(\sin \sqrt{R - A^2} \right)}{\det \left(\sqrt{R - A^2} \int_{-1}^1 e^{tA} dt \right)},$$

$$\sum_{n=1}^{\infty} \operatorname{tr}(R(A^2 + (\pi n)^2 I)^{-1}) =$$

$$\operatorname{tr}\left(\iiint_{0 \le \tau_1 \le \tau_2 \le t \le 1} e^{(\tau_2 - 2t)A} Re^{(\tau_2 - 2\tau_1)} d\tau_2 d\tau_1 dt \left(\int_0^1 e^{-2tA} dt \right)^{-1} \right).$$

The right-hand side of the determinant formula has an obvious meaning also in the case of a degenerate $R-A^2$. If $m=1,\ A=a,\ R=a^2+b^2$, we get:

$$\prod_{n=1}^{\infty} \left(1 - \frac{a^2 + b^2}{a^2 + (\pi n)^2} \right) = \frac{a \sin b}{b \sin a},$$

an interpolation between the classical Euler identity and its hyperbolic version. The trace identity is essentially simplified if the matrices R and A commute. In the commutative case we obtain:

$$\sum_{n=1}^{\infty} \operatorname{tr}(R(A^2 + (\pi n)^2 I)^{-1}) = \frac{1}{2} \operatorname{tr}\left(R(A \operatorname{cth} A - I)A^{-2}\right).$$

4 Proof of Theorem 1

We start with some definitions and notations. A compact quadratic form b on the Hilbert space \mathcal{V} is a form defined by a compact symmetric operator $B, b(v) = \langle Bv | v \rangle, v \in \mathcal{V}$. The spectrum of b and its positive and negative parts are, by the definition, those of B, i. e. $\mathrm{Sp}_{\pm}(b) = \mathrm{Sp}_{\pm}(B)$. Recall that the eigenvalues of B are just critical values of the restriction of b to the unit sphere in \mathcal{V} .

Let $\mathcal{V}_0 \subset \mathcal{V}$ be a Hilbert subspace, then the form $b|_{\mathcal{V}_0}$ is defined by the composition of $B|_{\mathcal{V}_0}$ and the orthogonal projection of \mathcal{V} on \mathcal{V}_0 . The relation between $\operatorname{Sp}(b)$ and $\operatorname{Sp}(b|_{\mathcal{V}_0})$ is ruled by the classical Rayleigh–Courant minimax principle (see [4, 5]).

Assume that $\operatorname{Sp}_{\pm}(b)$ is infinite; then a natural ordering β_n , $n \in \mathbb{Z}_{\pm}$, of $\operatorname{Sp}_{\pm}(b)$ is defined in the same way as the natural ordering of $\operatorname{Sp}_{\pm}(K)$ (see (7)).

Definition 3. We say that b has the spectrum of capacity $\varsigma > 0$ with the remainder of order $\nu > 1$ if $\mathrm{Sp}_{+}(b)$ and $\mathrm{Sp}_{-}(b)$ are both infinite and

$$\beta_n = \frac{\varsigma}{n} + O(n^{-\nu}) \quad \text{as} \quad n \to \pm \infty.$$
 (11)

We say that b has the spectrum of zero capacity with the remainder of order ν if either $\operatorname{Sp}_{\pm}(b)$ is finite or $\beta_n = O(n^{-\nu})$ as $n \to \pm \infty$.

Let b_i be a quadratic form on the Hilbert space \mathcal{V}_i , i = 1, 2; then $b_1 \oplus b_2$ is a quadratic form on $\mathcal{V}_1 \oplus \mathcal{V}_2$, $\operatorname{Sp}(b_1 \oplus b_2) = \operatorname{Sp}(b_1) \cup \operatorname{Sp}(b_2)$ and the multiplicities of common eigenvalues are added.

- **Proposition 3.** (i) If b has the spectrum of capacity $\varsigma \geq 0$, then sb has the spectrum of capacity $s\varsigma$ with the remainder of the same order as b, for any $s \in \mathbb{R}$.
- (ii) If b_1, b_2 have the spectra of capacities ς_1, ς_2 with the remainders of equal orders, then $b_1 \oplus b_2$ has the spectrum of capacity $\varsigma_1 + \varsigma_2$ and the remainder of the same order as b_1, b_2 .
- (iii) Let V_0 be a Hilbert subspace of the Hilbert space V and $\dim(V/V_0) < \infty$. Assume that one of two forms b or $b|_{V_0}$ has the spectrum of capacity $\varsigma \geq 0$ with a remainder of order $\nu \leq 2$. Then the second form has the spectrum of the the same capacity ς with a remainder of the same order ν .
- (iv) Let the forms b and \hat{b} be defined on the same Hilbert space \mathcal{V} , where b has the spectrum of capacity ς and \hat{b} has the spectrum of zero capacity, both with the reminder term of order ν . Then the form $b+\hat{b}$ has the spectrum of capacity ς with the reminder term of order $\frac{2\nu+1}{\nu+1}$.

Proof. Statement (i) is obvious. To prove (ii) we re-write asymptotic relation (11) in a more convenient form. An equivalent relation for positive n reads:

$$\#\left\{k \in \mathbb{Z} : 0 < \frac{1}{\beta_k} < n\right\} = \varsigma n + O(n^{2-\nu}), \text{ as } n \to \infty$$

and similarly for negative n. Statement (ii) follows immediately.

Statement (iii) follows from the Rayleigh–Courant minimax principle for the eigenvalues and the relation: $\left|\frac{\varsigma}{n} - \frac{\varsigma}{n+j}\right| = O(\frac{1}{n^2})$ as $|n| \to \infty$ for any fixed j.

To prove (iv) we use the Weyl inequality for the eigenvalues of the sum of two forms. Weyl inequality is a straightforward corollary of the minimax principle, it claims that the positive eigenvalue number i+j-1 in the natural ordering of the sum of two forms does not exceed the sum of the eigenvalue number i of the first summand and the eigenvalue number j of the second summand. Of course, we may equally works with naturally ordered negative eigenvalues simply changing the signs of the forms.

In our case, to have both sides estimates we first present $b+\hat{b}$ as the sum of b and \hat{b} and then present b as the sum of $b+\hat{b}$ and $-\hat{b}$. In the first case we apply the Weyl inequality with $i=n-[n^{\delta}],\ j=[n^{\delta}]$ for some $\delta\in(0,1)$, and in the second case we take $i=n,\ j=[n^{\delta}]$. The best result is obtained for $\delta=\frac{1}{\nu+1}$.

We have to prove that the spectrum of operator K (see (6), (5)) has capacity $\frac{1}{\pi} \int_0^1 \bar{\zeta}_t dt$ with the remainder of order $\frac{5}{3}$. First, we may identify $T_{\tilde{u}(t)}$ with \mathbb{R}^m and assume that $H_t = -I$. Indeed, if we trivialize the vector bundle $\tilde{u}^*(TU)$ over the segment [0,1], then H_t becomes a negative definite symmetric matrix. Then we substitute v by $(-H_t)^{\frac{1}{2}}v$ and Z_t by $Z_t(-H_t)^{-\frac{1}{2}}$.

Let $0 = t_0 < t_1 < \cdots < t_l < t_{l+1} = 1$ be a subdivision of the segment [0, 1]. The subspace

$$\{v \in \mathcal{V}: \int_{t_i}^{t_{i+1}} X_t v(t) dt = 0, \ i = 0, 1, \dots, l\} = \bigoplus_{i=0}^l \mathcal{V}_i, \ \mathcal{V}_i \subset L^2([0, 1]; \mathbb{R}^m)$$

has a finite codimension in \mathcal{V} . The quadratic form

$$\langle Kv|v \rangle = \int_0^1 \left\langle JZ_t v(t), \int_0^t Z_\tau v(\tau) d\tau \right\rangle dt$$
 (12)

restricted to this subspace turns into the direct sum of the forms

$$\langle K_i v | v \rangle = \int_{t_i}^{t_{i+1}} \left\langle J Z_t v(t), \int_{t_i}^t Z_\tau v(\tau) d\tau \right\rangle dt, \quad v_i \in \mathcal{V}_i, \ i = 0, 1 \dots l.$$

Indeed, the relations $\int_{t_i}^{t_{i+1}} X_t v(t) dt = 0$ imply that

$$\left\langle \int_{t_i}^{t_{i+1}} JZ_t v(t) dt, \int_{t_j}^{t_{j+1}} Z_t v(t) dt \right\rangle = 0.$$

According to Proposition 3 it is sufficient to prove our theorem for the operators K_i , i = 0, 1, ..., l. In particular, we may substitute the piecewise analyticity assumption in the statement of the theorem by the analyticity one.

Moreover, under the analyticity condition we may assume that

$$Z_t^* J Z_t = \bigoplus_{j=1}^k \begin{pmatrix} 0 & -\zeta_j(t) \\ \zeta_j(t) & 0 \end{pmatrix},$$

where $0 \le 2k \le m$ and $\zeta_j(t)$ are not identical zero. Indeed, according to the Rayleigh theorem (see [5]), there exists an analytically depending on t orthonormal basis in which our anti-symmetric matrix takes a desired form.

The functions $\zeta_j(t)$, $j=1,\ldots,k$, are analytic and may have only isolated zeros. Hence we may take a finer subdivision of [0,1] in such a way that $\zeta_j(t)$, $j=1,\ldots,k$, do not change sign on the segments $[t_i,t_{i+1}]$. Moreover, a simple change of the basis of \mathbb{R}^m if necessary allows us to assume that $\zeta_j(t) \geq 0$, $t \in [t_i,t_{i+1}]$. Actually, to simplify notations a little bit, we may simply assume that $\zeta_j(t) \geq 0$, $0 \leq t \leq 1$, $j=1,\ldots,k$. In this case $\bar{\zeta}(t) = \sum_{i=1}^k \zeta_j(t)$.

Let us study quadratic form (12) on the space $\{v \in L^2([0,1]; \mathbb{R}^m) : \int_0^1 v(t) dt = 0\}$. Recall that we are allowed by Proposition 3 to work on any subspace of $L^2([0,1]; \mathbb{R}^m)$ of a finite codimension. We set $w(t) = \int_0^t v(\tau) d\tau$; a double integration by parts gives:

$$\int_0^1 \left\langle J Z_t v(t), \int_0^t Z_\tau v(\tau) d\tau \right\rangle dt = \int_0^1 \left\langle J Z_t v(t), Z_t w(t) \right\rangle dt +$$

$$\int_0^1 \left\langle J Z_t w(t), \dot{Z}_t w(t) \right\rangle dt + \int_0^1 \left\langle J \dot{Z}_t w(t), \int_0^t \dot{Z}_\tau w(\tau) d\tau \right\rangle dt.$$

Moreover, we have:

$$\left| \int_0^1 \left\langle J Z_t w(t), \dot{Z}_t w(t) \right\rangle \, dt + \int_0^1 \left\langle J \dot{Z}_t w(t), \int_0^t \dot{Z}_\tau w(\tau) \, d\tau \right\rangle \, dt \right| \leq c \int_0^1 |w(t)|^2 \, dt$$

for some constant c.

Let $\bar{\lambda}_n$ and $\underline{\lambda}_n$, $n \in \mathbb{Z} \setminus \{0\}$, be naturally ordered non zero eigenvalues of the quadratic forms

$$\int_0^1 \langle JZ_t v(t), Z_t w(t) \rangle dt + c \int_0^1 |w(t)|^2 dt$$

and

$$\int_0^1 \langle JZ_t v(t), Z_t w(t) \rangle dt - c \int_0^1 |w(t)|^2 dt$$

correspondently. The minimax principle for the eigenvalues implies that $\underline{\lambda}_n \leq \alpha_n \leq \overline{\lambda}_n, \ n \in \mathbb{Z} \setminus \{0\}$. Moreover, the form

$$\int_0^1 \langle JZ_t v(t), Z_t w(t) \rangle dt \pm c \int_0^1 |w(t)|^2 dt$$

splits in the direct sum of the forms

$$\int_0^1 \zeta_j(t) \langle J v_j(t), w_j(t) \rangle dt \pm c \int_0^1 |w_j(t)|^2 dt, \quad j = 1, \dots, k,$$
 (13)

where $v_j(t) \in \mathbb{R}^2$, $w_j(t) = \int_0^t v_j(t)$, $w_j(1) = 0$, $J = \begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix}$ and we simply ignore the identically vanishing part if 2k < m.

It remains to estimate the spectrum of the forms (13). To do that, we study the spectrum of the forms

$$\int_0^1 \zeta_j(t) \langle Jv_j(t), w_j(t) \rangle dt, \qquad \int_0^1 |w_j(t)|^2 dt, \tag{14}$$

and apply the statement (iv) of Proposition 3 to their linear combinations. To simplify calculations, we identify \mathbb{R}^2 with \mathbb{C} as follows: $(v_j^1, v_j^2) = v_j^1 + i v_j^2$, then J is the multiplication on the imaginary unit i. A complex-valued function $v_j(\cdot)$ is an eigenfunction with an eigenvalue λ for the first of two forms (14) if and only if it is a critical point of the functional

$$\int_0^1 \langle \zeta(t)iv_j(t), w_j(t) \rangle - \lambda \langle v_j(t), v_j(t) \rangle dt \quad w_j(0) = w_j(1) = 0.$$

The Euler–Lagrange equation for this functional reads:

$$2\lambda \ddot{w}_i + (\zeta_i i w_i) + \zeta_i i \dot{w} = 0, \quad w_i(0) = w_i(1) = 0.$$

We set $x(t) = e^{\frac{i}{2\lambda} \int_0^t \zeta_j(\tau) d\tau} w_j(t)$, plug-in this expression in the equation and arrive to a standard Sturm-Liouville problem:

$$\ddot{x} + \left(\frac{\zeta_j}{2\lambda}\right)^2 x = 0, \quad x(0) = x(1) = 0.$$

We see that the spectrum is double (recall that $x \in \mathbb{C}$) and is symmetric with respect to the origin. The asymptotics of the spectrum for the Sturm–Liouville problem is well-known (see [4]). We obtain:

$$\lambda_n = \frac{1}{\pi n} \int_0^1 \zeta_j(t) dt + O\left(\frac{1}{n^2}\right), \quad n \to \pm \infty.$$

In other words, this spectrum has capacity $\frac{1}{\pi} \int_0^1 \zeta_j(t) dt$ and the remainder term of order 2. The Euler-Lagrange equation for the eigenfunctions problem for the second of two forms (14) reads: $\lambda \ddot{w}_j + w_j = 0$, $w_j(0) = w_j(1) = 0$; hence $\lambda_n = \frac{1}{(\pi n)^2} = O(\frac{1}{n^2})$. The spectrum of this form has zero capacity with the remainder of order 2.

5 Proof of Theorem 2

We again assume that $H_t = -I$ using the same preliminary change of variables as in the proof of Theorem 1, if necessary: substitute v by $(-H)^{\frac{1}{2}}v$ and Z_t by $Z_t(-H)^{-\frac{1}{2}}$. Moreover, we fix some coordinates in a neighborhood of $q_0 \in M$ and use induced coordinates in T^*M so that $T_{\tilde{\lambda}_0}(T^*M)$ is identified with $\mathbb{R}^d \times \mathbb{R}^d = \{(p,q): p,q \in \mathbb{R}^d\}$, where $d = \dim M$, $\pi_*(p,q) = q$, and

$$\sigma_{\tilde{\lambda}_0}\left((p,q),(p',q')\right) = \left\langle J\begin{pmatrix} p\\q \end{pmatrix},\begin{pmatrix} p'\\q' \end{pmatrix} \right\rangle, \quad J = \begin{pmatrix} 0 & -I\\I & 0 \end{pmatrix}.$$

The map $Z_t: \mathbb{R}^m \to \mathbb{R}^d \times \mathbb{R}^d$ has a form $Z_t: v \mapsto (Y_t v, X_t v)$, where Y_t and X_t are $d \times m$ -matrices. Let $s \in \mathbb{C}$; we define $Z_t^s: \mathbb{C}^d \times C^d$ by the formula $Z_t^s v = (sY_t v, X_t v)$. Now consider the complexified Jacobi equation:

$$\dot{\eta} = Z_t^s Z_t^{s*} J \eta, \quad \eta \in \mathbb{C}^d \times \mathbb{C}^d, \tag{15},$$

where the transposition "*" corresponds to the complex inner product and not to the Hermitian one! The matrix form of this equation is as follows:

$$\begin{pmatrix} \dot{p} \\ \dot{q} \end{pmatrix} = \begin{pmatrix} sY_t \\ X_t \end{pmatrix} \left(sY_t^*, X_t^* \right) \begin{pmatrix} -q \\ p \end{pmatrix}. \tag{16}$$

Proposition 4. Let $\Phi_t^s : \mathbb{C}^d \times \mathbb{C}^d \to \mathbb{C}^d \times \mathbb{C}^d$, $0 \le t \le 1$ be the fundamental matrix of the complexified Jacobi equation, i. e.

$$\frac{d}{dt}\Phi_t^s = Z_t^s Z_t^{s*} J \Phi_t^s, \quad \Phi_0^s = I;$$

then $s \mapsto \Phi_1^s$, $s \in \mathbb{C}$, is an entire matrix-valued function and

$$\|\Phi_1^s\| \le ce^{c|s|}$$

for some constant c

Proof. The Volterra series for our system has the form:

$$\Phi_1^s = I + \sum_{n=1}^{\infty} \int_0^1 \int_0^{t_1} \cdots \int_0^{t_{n-1}} Z_{t_1}^s Z_{t_1}^s J \cdots Z_{t_n}^s Z_{t_n}^s J dt_n \dots dt_1.$$

The *n*-th term of the series is a polynomial matrix of s whose norm is bounded by $\frac{1}{n!} \left(\int_0^1 \|Z_t^s Z_t^{s*}\| \, dt \right)^n$ for any s and n. The sum of the series is an entire function by the Weierstrass theorem. Moreover, we obtain the estimate: $\|\Phi_t^s\| \leq e^{\int_0^1 \|Z_t^s Z_t^{s*}\| \, dt}$ but it is worse than one we need.

To obtain a better estimate we make a change of variables in the matrix form of the complexified Jacobi equation (see (16)). We set $\eta^s=(p,sq)$; then $\dot{\eta}^s=sZ_tZ_t^*J\eta^s$ and we obtain: $|\eta^s(1)|\leq e^{|s|\int_0^1\|Z_tZ_t^*\|dt}|\eta^s(0)|$ that gives us the desired estimate of $\|\Phi_1^s\|$ for s separated from zero. \square

We define linear maps $Q_t^s: \mathbb{C}^n \to \mathbb{C}^n$ by the formula: $Q_t^s p(0) = q(t)$, where $t \mapsto \binom{p(t)}{q(t)}$ is the solution of system (16) with the initial condition q(0) = 0. Note that the matrix Q_t^s is real if $s \in \mathbb{R}$; moreover, $Q_t^1|_{\mathbb{R}^d} = Q_t$ (linear maps Q_t were defined just before the statement of Theorem 2). Proposition 4 implies that $\|Q_1^s\| \leq ce^{c|s|}$.

Proposition 5. Let $s \in \mathbb{R} \setminus \{0\}$; the matrix Q_1^s is degenerate if and only if $-\frac{1}{s} \in \operatorname{Sp}(K)$; moreover, dim ker Q_1^s equals the multiplicity of the eigenvalue $-\frac{1}{s}$.

Proof. Let $v \in \mathcal{V}$; it is easy to see that

$$\langle JZ_{\tau}^{s}v(\tau), Z_{t}^{s}v(t)\rangle = s \langle JZ_{\tau}v(\tau), Z_{t}v(t)\rangle.$$

Hence

$$<(I+sK)v|v> = \int_0^1 \langle v(t),v(t)\rangle dt - \int_0^1 \left\langle J \int_0^t Z_\tau^s v(\tau) d\tau, Z_t^s v(t) \right\rangle dt.$$

We keep symbol \mathcal{V} for the complexification of \mathcal{V} , i. e. $\mathcal{V} = \mathbb{C} \otimes \mathcal{V}$. A vector-function $v \in \mathcal{V}$ is an eigenvector of K with the eigenvalue $-\frac{1}{s}$ if and only if

$$\int_0^1 \left\langle v(t) - Z_t^{s*} J \int_0^t Z_t^s v(\tau) d\tau, w(t) \right\rangle dt = 0, \quad \forall w \in \mathcal{V}, \tag{17}$$

where
$$V = \{ w \in L_2([0,1]; \mathbb{C}^d) : \int_0^1 X_t w(t) dt = 0 \}.$$

Equation $\int_0^1 X_t w(t) dt = 0$ can be rewritten as follows: $\int_0^1 \langle J\nu, Z_t w(t) \rangle dt = 0$ for any $\nu \in \ker \pi_*$. In other words, $\mathcal{V}^{\perp} = \{t \mapsto Z_t^{s*} J\nu : \nu \in \ker \pi_*\}$. Hence relation (17) is equivalent to the existence of $\nu \in \ker \pi_*$ such that

$$v(t) = Z_t^{s*} J\left(\int_0^t Z_\tau^s v(\tau) d\tau + \nu\right).$$

Moreover, the vector ν is unique for given vector-function v since the Volterra equation $v(t) = Z_t^{s*} J \int_0^t Z_\tau^s v(\tau) d\tau$ has only zero solution. We set $\eta(t) = \int_0^t Z_\tau^s v(\tau) d\tau + \nu$ and obtain that a vector-function v satisfies relation (17) if and only if

$$\dot{\eta} = Z_t^s Z_t^{s*} J \eta, \quad 0 \le t \le 1, \quad \eta(0), \eta(1) \in \ker \pi_*.$$
 (18)

It follows that dim ker Q_1^s is equal to the multiplicity of the eigenvalue $-\frac{1}{s}$ of the operator K plus the dimension of the space of constant solutions of equation (18) that belong to ker π_* . Let $(\nu, 0) \in \mathbb{C}^d \times \mathbb{C}^d$ be such a solution. We plug-in it in (18) and get $X_t X_t^* \nu = 0$, $\forall t \in [0, 1]$. Hence $\nu = 0$.

Corollary 2. The equation det $Q_1^s = 0$, $s \in \mathbb{C}$, has only real roots.

Indeed, the operator K is symmetric and has only real eigenvalues. \square

Proposition 6. Let $s_0 \in \mathbb{R} \setminus \{0\}$ and $\det Q_1^{s_0} = 0$. Then the multiplicity of the root s_0 of the equation $\det Q_1^s = 0$ is equal to $\dim \ker Q_1^{s_0}$.

Proof. We may assume that $s \in \mathbb{R}$ and work in the real setting. We denote by $\eta^s(t) = (p^s(t), q^s(t)), \ t \in [0, 1]$, those solutions of the Hamiltonian system

$$\dot{\eta}^s = s Z_t Z_t^* J \eta^s \tag{17}$$

that satisfy the initial condition $q^s(0) = 0$; then $q^s(t) = sQ_t^sp^s(0)$ (c. f. proof of Proposition 4).

Let $\Lambda_1^s = \{(p^s(1), q^s(1)) : p^s(0) \in \mathbb{R}^d\}$; then Λ_1^s is a Lagrange subspace of the symplectic space $\mathbb{R}^d \times \mathbb{R}^d$ endowed with the standard symplectic structure $((p_1, q_1), (p_2, q_2)) \mapsto \langle p_1, q_2 \rangle - \langle p_2, q_1 \rangle$.

We take another Lagrange subspace Δ that is transversal to ker π_0^* and Λ_1^* . Then $\Delta = \{(Aq,q) : q \in \mathbb{R}^d\}$ where A is a symmetric matrix. We make a symplectic change of variables $(p,q) \mapsto (p',q)$ by putting p' = p - Aq. In new variables, Λ_1^s is transversal to the "horizontal" subspace defined by the equation p' = 0 for all s close to s_0 . Hence

$$\Lambda_1^s = \{ (p', R(s)p') : p' \in \mathbb{R}^d \},$$

where R(s) is a symmetric matrix. According to our construction, sQ_1^s is the product of an analytic with respect to s nondegenerate matrix and R(s). It follows that dim ker $R(s) = \dim \ker Q_1^s$ and the multiplicities of the root s_0 of the equations $\det R(s) = 0$ and $\det Q_1^s = 0$ are equal. We denote this multiplicity by $\mu(s_0)$.

Standard perturbation theory for the eigenvalues of symmetric operators (see, for instance, [5]) implies that: (i) $\mu(s_0) \geq \dim \ker R(s_0)$; (ii) if $\det \frac{dR}{ds}(s_0) \neq 0$, then $\mu(s_0) = \dim \ker R(s_0)$.

Now we have to compute $\frac{dR}{ds}$. To do that, we re-write equation (17) in coordinates (p',q). In new coordinates, this Hamiltonian equation has the same structure as (17) with $Z_t = (Y_t, X_t)$ substituted by $Z'_t = (Y_t + AX_t, X_t)$. From now on, we'll write p, Z_t instead of p', Z'_t and work directly with equation (17) in order to simplify notations. We have $\eta^s(1) = (p^s(1), R(s)q^s(1))$; the symmetricity of the matrix R(s) implies that

$$\left\langle J\eta^s(1), \frac{\partial \eta^s}{\partial s}(1) \right\rangle = \left\langle p^s(1), \left(\frac{dR}{ds}\right) p^s(1) \right\rangle.$$

Proposition 6 is now reduced to the following lemma:

Lemma 1. If
$$\left\langle J\eta^s(1), \frac{\partial \eta^s}{\partial s}(1) \right\rangle = 0$$
, then $\eta^s(1) = 0$.

Proof. We have: $\frac{\partial \dot{\eta}^s}{\partial s} = sZ_tZ_t^*J\frac{\partial \eta^s}{\partial s} + Z_tZ_t^*J\eta^s$. Hence

$$\frac{\partial}{\partial t} \left\langle J\eta^{s}(t), \frac{\partial \eta^{s}}{\partial s}(t) \right\rangle = s \left\langle JZ_{t}Z_{t}^{*}J\eta^{s}(t), \frac{\partial \eta^{s}}{\partial s}(t) \right\rangle + s \left\langle J\eta^{s}(t), Z_{t}Z_{t}^{*}J\frac{\partial \eta^{s}}{\partial s}(t) \right\rangle + \left\langle J\eta^{s}(t), Z_{t}Z_{t}^{*}J\eta^{s} \right\rangle.$$

First two terms in the right hand side of the last identity cancel because the matrix J is anti-symmetric and the matrix $JZ_tZ_t^*J$ is symmetric. We obtain:

$$\frac{\partial}{\partial t} \left\langle J \eta^s(t), \frac{\partial \eta^s}{\partial s}(t) \right\rangle = \left\langle Z_t^* J \eta^s(t), Z_t^* J \eta^s(t) \right\rangle.$$

Moreover, $\langle J\eta^s(0), \frac{\partial \eta^s}{\partial s}(0) \rangle = 0$ since $q^s(0) = 0$ for all s. We get:

$$\left\langle J\eta^s(1), \frac{\partial \eta^s}{\partial s}(1) \right\rangle = \int_0^1 |Z_t^* J\eta^s(t)|^2 dt.$$

If $\left\langle J\eta^s(1), \frac{\partial \eta^s}{\partial s}(1) \right\rangle = 0$, then $Z_t^*J\eta^s(t) = 0$, $\forall t \in [0,1]$. Equation (17) now implies that $\dot{\eta}^s(t) \equiv 0$ and $\eta^s(t) = (p^s(0),0)$. We plug-in this equilibrium in (17) and obtain: $X_tX_t^*p^s(0) = 0$, $0 \le t \le 1$. Hence $\Gamma_1p^s(0) = \int_0^1 X_tX_t^*dtp^s(0) = 0$. Recall that Γ_1 is a nondegenerate matrix. \square

Proposition 4 implies that $s\mapsto \det Q_1^s,\ s\in\mathbb{C}$, is an entire function that satisfies the estimate $\det Q_1^s\le ce^{c|s|}$ for some constant c. It follows from a classical Hadamard theorem (see [3]) that such a function has a presentation det $Q_1^s = ae^{bs}\prod_i(1-\frac{s}{s_i})$, where s_i are roots of the equation $\det Q_1^s = 0$ counted according to multiplicity and $a, b \in \mathbb{C}$.

Propositions 5 and 6 now imply that

$$\det Q_1^s = ae^{bs} \det(I + sK), \quad s \in \mathbb{C}.$$
(18)

It remains to find a and b. We have: $\det Q_1^0 = a$. Moreover,

$$\frac{d}{ds} \det Q_1^s|_{s=0} = \operatorname{tr} \left(\frac{dQ_1^s}{ds} \Big|_{s=0} (Q_1^0)^{-1} \right) \det Q_1^0.$$

We differentiate identity (18) and obtain: $\operatorname{tr}\left(\frac{dQ_1^s}{ds}\big|_{s=0}(Q_1^0)^{-1}\right) = b + \operatorname{tr} K$.

Let $\Pi:L^2([0,1];\mathbb{R}^m)\to \mathcal{V}$ be the orthogonal projection, then $\operatorname{tr} K=$ $\operatorname{tr}(\Pi \hat{K} \Pi) = \operatorname{tr}(K \Pi)$, where $\hat{K}v(t) = -\int_0^t Z_t^* J Z_\tau v(\tau) d\tau$. The subspace \mathcal{V} is the kernel of the operator $A: v \mapsto \int_0^1 X_t v(t) dt$, $v \in L^2([0,1]; \mathbb{R}^m)$. We have: $\Pi = I - A^*(AA^*)^{-1}A$, $\Pi v(t) = v(t) - X_t^* \Gamma_1^{-1} \int_0^1 X_\tau v(\tau) d\tau$.

The trace of a trace-class operator \bar{B} of a form $(\bar{B}v)(t) = \int_0^1 B(t,\tau)v(\tau) d\tau$ is computed according to the formula $\operatorname{tr} \bar{B} = \int_0^1 \operatorname{tr} B(t,t) \, dt$. This presentation of the trace is valid in our situation as well. Moreover, $\operatorname{tr}(Z_t^*JZ_t)\equiv 0$; we collect the terms and obtain the formula:

$$\operatorname{tr} K = \int_0^1 \int_0^t \operatorname{tr} \left(Z_t^* J Z_\tau X_\tau^* \Gamma_1^{-1} X_t \right) d\tau dt$$
$$= \int_0^1 \int_0^t \operatorname{tr} \left(X_t Z_t^* J Z_\tau X_\tau^* \Gamma_1^{-1} \right) d\tau dt$$

as in the statement of Theorem 2 (recall that we are working in coordinated where $H_t = -I$).

Now we compute Q_1^0 . For s=0, system (18) is reduced to the system: $\dot{p} = 0, \ \dot{q} = X_t X_t^* p$. Hence $Q_1^0 = \int_0^1 X_t X_t^* dt = \Gamma_1$. The last step is the calculation of $\frac{dQ_1^s}{ds}\big|_{s=0}$. System (16) reads:

$$\dot{p} = sY_t X_t^* p - s^2 Y_t Y_t^* q$$
$$\dot{q} = X_t X_t^* p - sX_t Y_t^* q.$$

We have:

$$q^{s}(t) = \int_{0}^{t} X_{\tau} X_{\tau}^{*} d\tau p_{0} + O(s), \quad p^{s}(t) = s \int_{0}^{t} Y_{\tau} X_{\tau}^{*} d\tau p_{0} + O(s^{2}).$$

Hence

$$q^{s}(1) = \int_{0}^{1} X_{t} X_{t}^{*} p(t)^{s} - s X_{t} Y_{t}^{*} q^{s}(t) dt =$$

$$\Gamma_{1} p_{0} + s \int_{0}^{1} \int_{0}^{t} X_{t} X_{t}^{*} Y_{\tau} X_{\tau}^{*} - X_{t} Y_{t}^{*} X_{\tau} X_{\tau}^{*} d\tau dt p_{0} + O(s^{2}).$$

Moreover, $X_t X_t^* Y_\tau X_\tau^* - X_t Y_t^* X_\tau X_\tau^* = X_t Z_t^* J Z_\tau X_\tau^*$. We see that $\operatorname{tr}\left(\frac{dQ_1^s}{ds}\big|_{s=0}\Gamma_1^{-1}\right) = \operatorname{tr} K$. Hence b=0.

Remark 5. We used only the first derivative of Q_1^s at s=0 in the proof of Theorem 2 but system (16) allows us to find explicit integral expressions for all higher derivatives and thus to obtain integral expressions for all elementary symmetric functions of the eigenvalues of the operator K.

References

- [1] A. Agrachev, Yu. Sachkov, Control Theory from the geometric viewpoint. Springer Verlag, 2004, 426 p.
- [2] S. Bolotin, D. Treschev, Hill's formula. Russian Math. Surveys, 2010, v. 65, 191–257
- [3] J. Conway, Functions on one complex variable, I. Springer Verlag, 1995, 317 p.
- [4] R. Courant, D. Hilbert, *Methods of mathematical physics*, *I.* Intersience Pub., 1953, 576 p.
- [5] T. Kato, Perturbation theory for linear operators. Springer Verlag, 1980, 643 p.