

Report No. 53/2004

Spectral Analysis of Partial Differential Equations

Organised by
Alexander V. Sobolev (Brighton)
Timo Weidl (Stuttgart)

November 28th – December 4th, 2004

ABSTRACT. In this workshop talks on discrete spectra, including Lieb-Thirring estimates and properties of resonances were presented. The connection between continuous and discrete spectra was discussed. Various characteristics of continuous spectra in the context of random or magnetic operators were investigated. The Liouville Theorem, the problem of absolute continuity, and the classical problem of homogenization for periodic operators were treated.

Mathematics Subject Classification (2000): 35xx, 46xx, 47xx, 81xx .

Introduction by the Organisers

The goal of the workshop was to bring together specialists working in various branches of spectral theory with applications to solid state physics, superconductivity, quantum mechanics etc. The meeting was attended by more than 45 participants from Europe, Japan, Russia, South America and US. During the five days 26 talks were delivered. A special care was taken that apart from the recognized experts in the field, young participants also had an opportunity to speak about their results. The Wednesday morning session, preceding the traditional afternoon hike, consisted of talks of survey nature, which was appreciated by all.

There were several major themes in the workshop. One was the study of discrete spectra, including Lieb-Thirring estimates, properties of resonances. A substantial number of talks was concerned with the connection between the continuous and discrete spectra. These include, in particular, the study of the so-called trace formulas. The investigation of various characteristics of the continuous spectra (e.g. density of states, spectral shift function) was featured in a number of talks in the context of random or magnetic operators. A variety of new results were also reported on the theory of periodic operators. They concerned the Liouville Theorem, the problem of absolute continuity, and the classical problem of homogenization.

A relatively low number of talks gave the participants an opportunity for discussions in small groups outside the scheduled lecture time. It is hoped that these contacts will result in further collaboration.

It is our pleasure to thank the administration and staff of the *Mathematisches Forschungsinstitut Oberwolfach* for creating comfortable and genuinely inspiring atmosphere, which facilitated the work of the organisers and contributed to the success of the workshop.

Workshop: Spectral Analysis of Partial Differential Equations

Table of Contents

Arne Jensen (joint with Gheorghe Nenciu)	
<i>A Fermi Golden Rule at thresholds</i>	2843
M. Sh. Birman (joint with T. A. Suslina)	
<i>Homogenization of periodic differential operators in \mathbb{R}^d as a spectral threshold effect</i>	2845
G.M. Graf (joint with F. Bernasconi, D. Hasler)	
<i>Finite Casimir energy for the electromagnetic field in a cavity</i>	2847
Frédéric Klopp (joint with Werner Kirsch)	
<i>The band-edge behavior of the density of surface states</i>	2850
Alexander Kiselev	
<i>Recent results on singular spectrum of Schrödinger operators</i>	2854
Kenji Yajima	
<i>Dispersive estimates for Schrödinger equations</i>	2857
Ira Herbst	
<i>Classical and Quantum Mechanics for a Particle in a Long-Range Magnetic Field</i>	2860
Jean-Marie Barbaroux (joint with V. Bach, M. Esteban, W. Farkas, B. Helffer, E. Séré and H. Siedentop)	
<i>Some variational principles for relativistic energy functionals</i>	2862
Mikhail Solomyak	
<i>On the mathematical model of the irreversible quantum graph</i>	2865
Rafael D. Benguria (joint with Michael Loss)	
<i>A Lieb–Thirring Inequality and an Isoperimetric Problem for Closed Curves Curves in \mathbb{R}^2</i>	2867
B. Helffer (joint with S. Fournais)	
<i>Accurate estimates for magnetic bottles in connection with superconductivity</i>	2869
Shu Nakamura	
<i>Wave Front Set for Solutions to Schrödinger Equations</i>	2872
D. R. Yafaev	
<i>Spectral shift function for self-adjoint operators without spectral gaps</i>	2875
László Erdős (joint with Jan Philip Solovej)	
<i>Uniform Magnetic Lieb–Thirring inequalities</i>	2876

Werner Kirsch	
<i>Old and New Tales about Lifshitz Tails</i>	2879
Peter Kuchment (joint with Yehuda Pinchover, Technion, Israel)	
<i>Liouville theorems on abelian coverings of compact manifolds</i>	2882
T. A. Suslina	
<i>Homogenization problem for the stationary periodic Maxwell system</i>	2885
David Damanik	
<i>Bound States and Essential Spectrum</i>	2888
Maria Hoffmann-Ostenhof (joint with Søren Fournais, Thomas Hoffmann-Ostenhof and Thomas Østergaard Sørensen)	
<i>Properties of Coulombic wavefunctions and their electron density</i>	2890
Hynek Kovařík (joint with Denis Borisov and Tomas Ekholm)	
<i>On a Magnetic Hardy Inequality in The Waveguide</i>	2891
Alexander Pushnitski (joint with Nikolai Filonov)	
<i>Spectral asymptotics for the Landau Hamiltonian and logarithmic capacity</i>	2893
Dirk Hundertmark (joint with Vadim Zharnitsky)	
<i>Gaussian extremizers for the Strichartz inequality in one and two dimensions</i>	2895
Daniel M. Elton	
<i>Some Results on the Spectra of Periodic Landau Operators</i>	2897
Georgi Raikov	
<i>Spectral Shift Function for Magnetic Schrödinger Operators</i>	2900
Rupert L. Frank	
<i>On the Laplacian in the halfspace with a periodic boundary condition</i>	2903
A. Laptev (joint with S. Naboko and O. Safronov)	
<i>A Multidimensional Trace Formula</i>	2905

Abstracts

A Fermi Golden Rule at thresholds

ARNE JENSEN

(joint work with Gheorghe Nenciu)

We describe our main results in the form of an example. Consider a Schrödinger operator

$$H = -\Delta + V \quad \text{on } L^2(\mathbf{R}^3),$$

where we assume that $V \in C_0^\infty(\mathbf{R}^3)$. The essential spectrum is $[0, \infty)$, and is purely absolutely continuous. There may be a finite number of negative eigenvalues, and an eigenvalue at zero. We assume here that 0 is a non-degenerate eigenvalue, with normalized eigenfunction Ψ_0 . We study what happens to this eigenvalue under small perturbations. Let $W \in C_0^\infty(\mathbf{R}^3)$. To avoid the case that 0 becomes a negative discrete eigenvalue, we introduce

Assumption (A1). $b = \langle \Psi_0, W\Psi_0 \rangle > 0$.

The we consider

$$H(\varepsilon) = H + \varepsilon W, \quad \varepsilon > 0.$$

We show that the zero eigenvalue becomes a resonance, in the time-dependent sense introduced by A. Orth [5]. In order to formulate the main result we need some further results and assumptions.

In the resolvent $R(z) = (H - z)^{-1}$ we change the variable to $\kappa = -i\sqrt{z}$, $\text{Im } z > 0$, $\text{Re } \kappa \geq 0$. It is well-known (see [2]) that we have an asymptotic expansion as $\kappa \rightarrow 0$

$$R(-\kappa^2) = \frac{1}{\kappa^2} P_0 + \sum_{j=-1}^N \kappa^j G_j + \mathcal{O}(\kappa^{N+1}),$$

valid in the topology of the weighted spaces, $\mathcal{B}(L^s(\mathbf{R}^3), L^{-s}(\mathbf{R}^3))$, for s sufficiently large, depending on N .

We can now formulate the next essential assumption.

Assumption (A2). *There exists an odd integer $\nu \geq -1$, such that*

$$g_\nu = \langle \Psi_0, W G_\nu W \Psi_0 \rangle \neq 0, \quad G_j = 0, \quad j = -1, 1, 3, \dots, \nu - 2.$$

Our main result can then be formulated as follows.

Theorem. *There exists $\varepsilon_0 > 0$ such that*

$$\langle \Psi_0, e^{-itH(\varepsilon)} \Psi_0 \rangle = e^{-it\lambda(\varepsilon)} + \delta(\varepsilon, t), \quad t > 0, \quad 0 < \varepsilon < \varepsilon_0.$$

Here

$$|\delta(\varepsilon, t)| \leq C\varepsilon^{p(\nu)} |\ln \varepsilon|^t,$$

where $\iota = 1$ for $\nu = -1, 1$, and zero otherwise. We write $p(\nu) = \min\{2, (2 + \nu)/2\}$. We have $\lambda(\varepsilon) = x_0(\varepsilon) - i\Gamma(\varepsilon)$, with the expansions

$$\begin{aligned}x_0(\varepsilon) &= b\varepsilon(1 + \mathcal{O}(\varepsilon)), \\ \Gamma(\varepsilon) &= -i^{\nu-1}g_\nu b^{\nu/2} \varepsilon^{2+(\nu/2)}(1 + \mathcal{O}(\varepsilon)),\end{aligned}$$

as $\varepsilon \rightarrow 0$.

We note that $-i^{\nu-1}g_\nu > 0$. Our main result holds in an abstract setting, where we assume the existence of an asymptotic expansion of the type above for the resolvent of H .

The result shows how the Fermi Golden Rule has to be modified, to get the lifetime of the resonance. Notice that in the usual case of perturbation of an eigenvalue embedded in the continuum proper, the coupling constant dependence for the imaginary part is ε^2 , whereas we have $\varepsilon^{2+(\nu/2)}$, $\nu \geq -1$ and odd. All possible values of ν can be shown to occur in explicit examples.

It is possible to compute g_ν explicitly. In the example under consideration we have the following result. Take Ψ_0 real-valued, and let

$$X_j = \int_{\mathbf{R}^3} \Psi_0(x)V(x)x_j dx, \quad j = 1, 2, 3.$$

Assume that at least one $X_j \neq 0$. Then $\nu = -1$, and we have

$$g_{-1} = \frac{b^2}{12\pi}(X_1^2 + X_2^2 + X_3^2).$$

Resonances can also be defined as poles of a meromorphic continuation of the resolvent $R(z)$ in a suitable sense. The problem of perturbation of a threshold eigenvalue was studied by B. Baumgartner [1] in a two channel setting, using meromorphic continuation. He also gave heuristics for the modification of the Fermi Golden Rule, which agrees with our main theorem.

In [4] we give complete results, and, based on the resolvent expansions in [3], we give a large number of examples of both one channel and two channel Schrödinger operators satisfying our assumptions.

REFERENCES

- [1] B. Baumgartner, *Interchannel resonances at a threshold*, J. Math. Phys. **37** (1996), 5928–5938.
- [2] A. Jensen and T. Kato, *Spectral properties of Schrödinger operators and time-decay of the wave functions*, Duke Math. J. **46** (1979), 583–611.
- [3] A. Jensen and G. Nenciu, *A unified approach to resolvent expansions at thresholds*, Rev. Math. Phys. **13** (2001), no. 6, 717–754.
- [4] A. Jensen and G. Nenciu, *The Fermi Golden Rule at thresholds*, in preparation.
- [5] A. Orth, *Quantum mechanical resonance and limiting absorption: the many body problem*, Comm. Math. Phys. **126** (1990), no. 3, 559–573.

Homogenization of periodic differential operators in \mathbb{R}^d as a spectral threshold effect

M. SH. BIRMAN

(joint work with T. A. Suslina)

In [1], the *spectral approach* to homogenization problems for one class of selfadjoint elliptic matrix second order differential operators is systematically developed. On the basis of the Floquet-Bloch decomposition, it is shown that homogenization is a *threshold effect* near the bottom of the spectrum. In what follows, we explain this point of view and discuss the typical results. The results presented in Section 4 are new.

1. Let $\Gamma \subset \mathbb{R}^d$ be a lattice, and let Ω be the cell of Γ . We use the notation $\mathfrak{G} = L_2(\mathbb{R}^d; \mathbb{C}^n)$, $\mathfrak{G}_* = L_2(\mathbb{R}^d; \mathbb{C}^m)$, $\mathbf{D} = -i\nabla$. It is assumed that $m \geq n$. Let h be an $(m \times m)$ -matrix-valued Γ -periodic function in \mathbb{R}^d such that $h, h^{-1} \in L_\infty(\mathbb{R}^d)$. We put $g = h^*h$. Let $b(\boldsymbol{\xi})$, $\boldsymbol{\xi} \in \mathbb{R}^d$, be an $(m \times n)$ -matrix-valued linear homogeneous function such that $\text{rank } b(\boldsymbol{\xi}) = n$ for $\boldsymbol{\xi} \neq 0$. Then

$$\alpha_0 \mathbf{1}_n \leq b(\boldsymbol{\theta})^* b(\boldsymbol{\theta}) \leq \alpha_1 \mathbf{1}_n, \quad |\boldsymbol{\theta}| = 1, \quad 0 < \alpha_0 \leq \alpha_1 < \infty.$$

We consider the first order differential operator $hb(\mathbf{D}) = \mathcal{X} : \mathfrak{G} \rightarrow \mathfrak{G}_*$; $\text{Dom } \mathcal{X} = H^1(\mathbb{R}^d; \mathbb{C}^n)$. Here H^1 is the Sobolev space. Then the operator $\mathcal{A}(g) = \mathcal{X}^* \mathcal{X} = b(\mathbf{D})^* g(\mathbf{x}) b(\mathbf{D})$ is selfadjoint in \mathfrak{G} .

Our main object is the operator family $\mathcal{A}_\varepsilon(g) = \mathcal{A}(g^\varepsilon)$, where $g^\varepsilon(\mathbf{x}) = g(\varepsilon^{-1}\mathbf{x})$, $\varepsilon > 0$. We study the behavior of solutions of the equation

$$\mathcal{A}_\varepsilon(g) \mathbf{u}_\varepsilon + \mathbf{u}_\varepsilon = \mathbf{F}, \quad \mathbf{F} \in \mathfrak{G}, \quad (1)$$

as $\varepsilon \rightarrow 0$.

2. The following definition of the constant effective matrix g^0 is standard for the homogenization theory. Let $\mathbf{C} \in \mathbb{C}^m$, and let \mathbf{w} be a weak Γ -periodic solution of the equation

$$b(\mathbf{D})^* g(\mathbf{x}) (b(\mathbf{D}) \mathbf{w} + \mathbf{C}) = 0. \quad (2)$$

Then g^0 is defined by the relation

$$g^0 \mathbf{C} = |\Omega|^{-1} \int_{\Omega} g(\mathbf{x}) (b(\mathbf{D}) \mathbf{w} + \mathbf{C}) d\mathbf{x}.$$

The effective matrix satisfies the estimates

$$|\Omega| \left(\int_{\Omega} (g(\mathbf{x}))^{-1} d\mathbf{x} \right)^{-1} = \underline{g} \leq g^0 \leq \bar{g} = |\Omega|^{-1} \int_{\Omega} g(\mathbf{x}) d\mathbf{x}.$$

Theorem 1. *We have*

$$\|(\mathcal{A}(g) + \varepsilon^2 I)^{-1} - (\mathcal{A}(g^0) + \varepsilon^2 I)^{-1}\|_{\mathfrak{G} \rightarrow \mathfrak{G}} \leq C \varepsilon^{-1}, \quad 0 < \varepsilon \leq 1, \quad (3)$$

where the constant C depends only on Γ , α_0 , α_1 , $\|h\|_{L_\infty}$, $\|h^{-1}\|_{L_\infty}$.

The estimate (3) is of *threshold nature*, since we consider the resolvent in point $(-\varepsilon^2)$, i. e., near the bottom of the spectrum. (We have $\inf \text{spec } \mathcal{A}(g) = 0$.)

3. Along with equation (1), we consider the *homogenized* equation $\mathcal{A}(g^0)\mathbf{u}_0 + \mathbf{u}_0 = \mathbf{F}$. By traditional means of homogenization theory, it is easily proved that \mathbf{u}_ε tends to \mathbf{u}_0 weakly in $H^1(\mathbb{R}^d; \mathbb{C}^n)$. Using the spectral approach, we prove the following result, which complements this statement essentially.

Theorem 2. *Let C be the constant from inequality (3). Then*

$$\|(\mathcal{A}_\varepsilon(g) + I)^{-1} - (\mathcal{A}(g^0) + I)^{-1}\|_{\mathfrak{G} \rightarrow \mathfrak{G}} \leq C\varepsilon, \quad 0 < \varepsilon \leq 1. \quad (4)$$

Apparently, estimates of the form (4) are new for homogenization theory. In fact, estimates (3) and (4) are equivalent. Indeed, let T_ε be the unitary scale transformation in \mathfrak{G} : $(T_\varepsilon \mathbf{u})(\mathbf{x}) = \varepsilon^{d/2} \mathbf{u}(\varepsilon \mathbf{x})$. Then

$$(\mathcal{A}_\varepsilon(g) + I)^{-1} = \varepsilon^2 T_\varepsilon^* (\mathcal{A}(g) + \varepsilon^2 I)^{-1} T_\varepsilon.$$

The operator $\mathcal{A}(g^0)$ satisfies similar identity, but $(g^0)^\varepsilon = g^0$.

Note that using the scale transformation is possible only for the estimates in the operator norm. For the study of convergence of different types, this method does not work.

4. In homogenization theory, adding appropriate *correction term* of order ε to \mathbf{u}_0 , one obtains more accurate approximation for \mathbf{u}_ε . This correction term contains some rapidly oscillating factors. This way is also possible in L_2 -theory. Here we present the corresponding result for the simplest case, namely, for the operator $\mathcal{A}(g) = \mathbf{D}^* g(\mathbf{x}) \mathbf{D} = -\text{div } g(\mathbf{x}) \nabla$ (now $n = 1$, $m = d$, $b(\boldsymbol{\xi}) = \boldsymbol{\xi}$). Let $v_j(\mathbf{x})$, $j = 1, \dots, d$, be the Γ -periodic solution of the equation $\mathbf{D}^* g(\mathbf{x})(\mathbf{D} v_j(\mathbf{x}) + \mathbf{e}_j) = 0$ such that $\int_\Omega v_j(\mathbf{x}) d\mathbf{x} = 0$. Here $\{\mathbf{e}_j\}$, $j = 1, \dots, d$, is the standard basis in \mathbb{R}^d . By $\Lambda(\mathbf{x})$ we denote the matrix-row $\{v_1(\mathbf{x}), v_2(\mathbf{x}), \dots, v_d(\mathbf{x})\}$. Then $\Lambda(\mathbf{x})$ is Γ -periodic. We put $\Lambda^\varepsilon(\mathbf{x}) = \Lambda(\varepsilon^{-1} \mathbf{x})$, and consider the operator $Z_\varepsilon : \mathfrak{G} \rightarrow \mathfrak{G}$,

$$Z_\varepsilon = \Lambda^\varepsilon \mathbf{D} (\mathcal{A}(g^0) + I)^{-1}.$$

Theorem 3. *For the operator $\mathcal{A}(g) = \mathbf{D}^* g(\mathbf{x}) \mathbf{D}$ under the above assumptions we have*

$$\|(\mathcal{A}_\varepsilon(g) + I)^{-1} - (\mathcal{A}(g^0) + I)^{-1} - \varepsilon(Z_\varepsilon + Z_\varepsilon^*)\|_{\mathfrak{G} \rightarrow \mathfrak{G}} \leq C_* \varepsilon^2, \quad 0 < \varepsilon \leq 1, \quad (5)$$

where C_* depends only on Γ , $\|g\|_{L^\infty}$, $\|g^{-1}\|_{L^\infty}$.

Remarks. 1) In homogenization theory, the traditional correction term is Z_ε . However, in L_2 -theory, in order to obtain the precise estimate (5), we have to take the symmetric expression $(Z_\varepsilon + Z_\varepsilon^*)$. In the case where the columns of $g(\mathbf{x})$ are divergence free, we have $Z_\varepsilon = 0$ and then

$$\|(\mathcal{A}_\varepsilon(g) + I)^{-1} - (\mathcal{A}(g^0) + I)^{-1}\|_{\mathfrak{G} \rightarrow \mathfrak{G}} \leq C_* \varepsilon^2, \quad 0 < \varepsilon \leq 1.$$

2) The estimate similar to (5) is true for the matrix operators $\mathcal{A}(g)$ defined in Section 1, if $m = n$. Apparently, in the general case one more additional summand should be added in the correction term.

REFERENCES

- [1] M. Sh. Birman and T. A. Suslina, *Second order periodic differential operators. Threshold properties and homogenization.*, Algebra i Analiz **15** (2003), 1–108; English transl., St. Petersburg Math. J. **15** (2004), 1–77.

Finite Casimir energy for the electromagnetic field in a cavity

G.M. GRAF

(joint work with F. Bernasconi, D. Hasler)

We present a Hilbert space formulation [3] of the *classical* Maxwell equations in a cavity $\Omega \subset \mathbb{R}^3$. In a preliminary Hilbert space $L^2(\Omega, \mathbb{C}^3)$ of (complex-valued) vector fields on Ω we define the dense subspaces

$$\begin{aligned} \mathcal{R} &= \{ \mathbf{V} \in L^2(\Omega, \mathbb{C}^3) \mid \operatorname{rot} \mathbf{V} \in L^2(\Omega, \mathbb{R}^3) \}, \\ \mathcal{R}_0 &= \{ \mathbf{V} \in \mathcal{R} \mid \langle \mathbf{U}, \operatorname{rot} \mathbf{V} \rangle = \langle \operatorname{rot} \mathbf{U}, \mathbf{V} \rangle, \forall \mathbf{U} \in \mathcal{R} \} \end{aligned}$$

and the (closed) operator $R = \operatorname{rot}$ with domain $\mathcal{D}(R) = \mathcal{R}_0$. Its adjoint is $R^* = \operatorname{rot}$ with $\mathcal{D}(R^*) = \mathcal{R}$. We remark that R , resp. R^* , is also the closure of rot defined on smooth vector fields \mathbf{V} with boundary condition $\mathbf{V}_{\parallel} = 0$ on the smooth boundary $\partial\Omega$, resp. without boundary conditions. Similarly, gradients $\nabla, \tilde{\nabla} : L^2(\Omega) \rightarrow L^2(\Omega, \mathbb{C}^3)$ can be defined with domains $\mathcal{D}(\nabla) = \{ \varphi \in L^2(\Omega) \mid \nabla\varphi \in L^2(\Omega), \varphi = 0 \text{ on } \partial\Omega \}$, resp. $\mathcal{D}(\tilde{\nabla})$ without the last boundary condition. Clearly, $\operatorname{Ran} \nabla \subset \operatorname{Ker} R$, $\operatorname{Ran} \tilde{\nabla} \subset \operatorname{Ker} R^*$, so that

$$(1) \quad \begin{aligned} \operatorname{Ran} R^* &\subset (\operatorname{Ker} R)^{\perp} \subset (\operatorname{Ran} \nabla)^{\perp} =: \mathcal{H}, \\ \operatorname{Ran} R &\subset (\operatorname{Ker} R^*)^{\perp} \subset (\operatorname{Ran} \tilde{\nabla})^{\perp} =: \mathcal{H}'. \end{aligned}$$

Therefore the Maxwell operator

$$M = \begin{pmatrix} 0 & iR^* \\ -iR & 0 \end{pmatrix} = M^*$$

on $L^2(\Omega, \mathbb{C}^3) \oplus L^2(\Omega, \mathbb{C}^3)$ restricts to the invariant subspace $\mathcal{H} \oplus \mathcal{H}'$, which is the *physical* Hilbert space for electromagnetic fields (\mathbf{E}, \mathbf{B}) . Indeed, the spaces

$$(2) \quad \begin{aligned} \mathcal{H} &= \{ \mathbf{E} \in L^2(\Omega, \mathbb{C}^3) \mid \operatorname{div} \mathbf{E} = 0 \}, \\ \mathcal{H}' &= \{ \mathbf{B} \in L^2(\Omega, \mathbb{C}^3) \mid \operatorname{div} \mathbf{B} = 0, \mathbf{B}_{\perp} = 0 \text{ on } \partial\Omega \} \end{aligned}$$

consist of divergence free fields and the Maxwell equations can be written as

$$(3) \quad i \frac{\partial}{\partial t} \begin{pmatrix} \mathbf{E} \\ \mathbf{B} \end{pmatrix} = M \begin{pmatrix} \mathbf{E} \\ \mathbf{B} \end{pmatrix}.$$

The usual boundary conditions $\mathbf{E}_{\parallel} = 0, \mathbf{B}_{\perp} = 0$ on the ideally conducting shell $\partial\Omega$ are accounted for through $\mathcal{D}(R)$, resp. \mathcal{H}' .

Remark. In [2] we defined M as an operator on $\mathcal{H} \oplus \mathcal{H}$. The difference consists of fields $(\mathbf{E}, \mathbf{B}) = (0, \nabla\psi)$ with ψ harmonic, and hence of (infinitely many) zero modes of M , which are irrelevant to the Casimir energy, see below.

We discuss the heat kernel traces for $M^2 = \text{diag}(R^*R, RR^*)$,

$$(4) \quad \text{Tr}_{\mathcal{H}}(e^{-tM^2}) = \sum_k e^{-t\omega_k^2} \cong \sum_{n=0}^{\infty} a_n t^{\frac{n-3}{2}}, \quad (t \downarrow 0),$$

(and similarly for \mathcal{H}' with coefficients a'_n), where ω_k^2 are the eigenvalues of R^*R on \mathcal{H} , resp. RR^* on \mathcal{H}' . They come in pairs, except for zero modes, and correspond to a single oscillator mode $\omega_k > 0$ for (3). The coefficients a_n are known, see e.g. [6, 4], for general operators of Laplace type. The direct application of such results is prevented by the divergence constraint in \mathcal{H} and \mathcal{H}' , see (2).

Let $L_{ab} = (\nabla_{\mathbf{e}_a} \mathbf{e}_b, \mathbf{n})$, ($a, b = 1, 2$), be the second fundamental form on the boundary $\partial\Omega$ with inward normal \mathbf{n} and local orthonormal frame $\{\mathbf{e}_1, \mathbf{e}_2, \mathbf{n}\}$. We denote by $|\Omega|$ the volume of Ω and set $f[\partial\Omega] = \int_{\partial\Omega} f(y) dy$, where dy is the (induced) Euclidean surface element on $\partial\Omega$. The corresponding Laplacian on $\partial\Omega$ is denoted by ∇^2 .

Theorem. [2] *Let $\Omega \subset \mathbb{R}^3$ be an open, connected domain with compact closure and smooth boundary $\partial\Omega$. Then*

$$\begin{aligned} a_0 &= 2(4\pi)^{-\frac{3}{2}} |\Omega|, & a_1 &= 0, & a_2 &= -\frac{4}{3}(4\pi)^{-\frac{3}{2}} (\text{tr}L)[\partial\Omega], \\ a_3 &= \frac{1}{64}(4\pi)^{-1} (3(\text{tr}L)^2 + 28 \det L)[\partial\Omega], \\ a_4 &= \frac{16}{315}(4\pi)^{-\frac{3}{2}} (2(\text{tr}L)^3 - 9\text{tr}L \cdot \det L)[\partial\Omega], \\ a_5 &= \frac{1}{122880}(4\pi)^{-1} (2295(\text{tr}L)^4 - 12440(\text{tr}L)^2 \det L + \\ &\quad + 13424(\det L)^2 + 1200\text{tr}L \cdot \nabla^2 \text{tr}L)[\partial\Omega]. \end{aligned}$$

The coefficients a'_n are the same, except for $n = 3$, where

$$a'_3 = \frac{1}{64}(4\pi)^{-1} (3(\text{tr}L)^2 - 36 \det L)[\partial\Omega] + 1.$$

By the Gauss-Bonnet theorem we have $a_3 - a'_3 = (4\pi)^{-1}(\det L)[\partial\Omega] - 1 = \sum_{i=1}^n (1 - g_i) - 1 = (n - 1) - \sum_{i=1}^n g_i$, where g_1, g_2, \dots, g_n are the genera of the n connected components of $\partial\Omega$. This equals the difference in the numbers of electrostatic, $n - 1$, and magnetostatic, $\sum_{i=1}^n g_i$, modes.

Sketch of proof. The transversal modes of the electromagnetic field, together with their unphysical, longitudinal counterparts in $\text{Ran } \nabla$ and $\text{Ran } \tilde{\nabla}$, see (1), are the eigenfunctions of the Laplacian acting on unconstrained vector fields, to which existing heat kernel expansions may be applied. The spurious contribution so introduced is essentially that of the Laplacian on scalar fields. Alternatively, consider the Laplacian of the de Rham complex of a 3-manifold with boundary. The electric and magnetic fields are then associated to forms of degree $p = 1$ and $p = 2$ respectively. In this correspondence transverse modes are associated with coexact, resp. exact forms, which permit to map longitudinal modes to forms of degree $p = 0$ and $p = 3$.

We apply the Theorem to the Casimir effect of the *quantum* field. To this end we retain: (i) The Weyl term a_0 is proportional to the volume of the cavity; (ii) $a_1 = 0$; (iii) a_2, a_4 are odd in the second fundamental form of the boundary; (iv) the asymptotic series (4) may be differentiated w.r.t. t . For the purpose of this discussion we simply define the Casimir energy by the mode summation method, see e.g. [1]. We shall observe that it is finite – a conclusion drawn in [1], but questioned in [9]. We do not however address the issue [7] of whether this definition is the most appropriate physically, nor do we compare it with others, based e.g. on the local energy density.

We enclose the cavity $\Omega \subset \mathbb{R}^3$ in a large ball Ω_0 and compare the vacuum energy of the electromagnetic field in the domains $\Omega \cup (\Omega_0 \setminus \overline{\Omega})$ with that of the reference domain Ω_0 . Each eigenmode of either configuration contributes a zero-point energy $\omega_k/2$, resp. $\omega_k^0/2$. As a regulator for the eigenfrequencies $\omega_k = \lambda_k^{1/2}$, we choose $e^{-\gamma\lambda_k}$, ($\gamma > 0$). The corresponding definition of the Casimir energy is

$$E_C = \frac{1}{2} \lim_{\Omega_0 \uparrow \mathbb{R}^3} \lim_{\gamma \downarrow 0} \left(\sum_k \lambda_k^{\frac{1}{2}} e^{-\gamma\lambda_k} - \sum_k (\lambda_k^0)^{\frac{1}{2}} e^{-\gamma\lambda_k^0} \right).$$

We now show that the limit $\gamma \downarrow 0$ is finite. (The subsequent limit $\Omega_0 \uparrow \mathbb{R}^3$ also exists.) Using $\lambda_k^{1/2} = -\pi^{-1/2} \int_0^\infty dt t^{-1/2} d(e^{-t\lambda_k})/dt$ we obtain

$$\begin{aligned} \sum_k \lambda_k^{\frac{1}{2}} e^{-\gamma\lambda_k} &\approx - \sum_{n=0}^4 \frac{n-3}{2\sqrt{\pi}} a_n \int_0^\delta dt t^{-\frac{1}{2}} (t+\gamma)^{\frac{n-5}{2}} \\ &\approx \frac{2}{\sqrt{\pi}} a_0 \gamma^{-2} + \frac{\sqrt{\pi}}{2} a_1 \gamma^{-\frac{3}{2}} + \frac{1}{\sqrt{\pi}} a_2 \gamma^{-1} + 0 \cdot a_3 \gamma^{-\frac{1}{2}} + \frac{1}{2\sqrt{\pi}} a_4 \log \gamma, \end{aligned}$$

where $\delta > 0$ is arbitrary, but fixed, and “ \approx ” means up to bounded terms as $\gamma \downarrow 0$. Hence a finite E_C requires that a_0, a_1, a_2, a_4 (but not necessarily a_3 !) agree for $\Omega \cup (\Omega_0 \setminus \overline{\Omega})$ and for the reference domain Ω_0 [8], [5]. By the Theorem this is so for a_0 and a_1 , but also for a_2, a_4 as the contribution from the two sides of $\partial\Omega$ cancel.

Acknowledgments. We thank M. Birman for suggesting the change in formulation mentioned in the remark, and him and K. Milton for pointing out to us refs. [3, 5], respectively.

REFERENCES

- [1] R. Balian, B. Duplantier, *Electromagnetic waves near perfect conductors, II. Casimir effect*, Ann. Phys. (N.Y.) **112** (1978), 165.
- [2] F. Bernasconi, G.M. Graf, D. Hasler, *The heat kernel expansion for the electromagnetic field in a cavity*, Ann. H. Poincaré **4** (2003), 1001.
- [3] M.Sh. Birman, M.Z. Solomyak, *The selfadjoint Maxwell operator in arbitrary domains* (Russian), Algebra i Analiz **1** (1989), 96; translation in Leningrad Math. J. **1** (1990), 99.
- [4] N. Blažić, N. Bokan, P.B. Gilkey, *Spectral geometry of the form valued Laplacian for manifolds with boundary*, Indian J. Pure Appl. Math. **23** (1992), 103.
- [5] M. Bordag, K. Kirsten, D. Vassilevich, *Ground state energy for a penetrable sphere and for a dielectric ball*, Phys. Rev. **D59** (1999), 085011.
- [6] T.P. Branson, P.B. Gilkey, K. Kirsten, D.V. Vassilevich, *Heat kernel asymptotics with mixed boundary conditions*, Nucl. Phys. **B563** (1999), 603.

- [7] P. Candelas, *Vacuum energy in the presence of dielectric and conducting surfaces*, Ann. Phys. (N.Y.) **143** (1982), 241.
- [8] G. Cognola, L. Vanzo, S. Zerbini, *Regularization dependence of vacuum energy in arbitrarily shaped cavities*, J. Math. Phys. **33** (1992), 222.
- [9] D. Deutsch, P. Candelas, *Boundary effects in quantum field theory*, Phys. Rev. D **20** (1978), 895.

The band-edge behavior of the density of surface states

FRÉDÉRIC KLOPP

(joint work with Werner Kirsch)

This talk is devoted to the integrated density of surface states for a simple discrete model of surface random operators (see e.g. [3, 1, 2, 7]). We study the asymptotic behavior of this quantity near the edges of the spectrum of the random model. The results are taken from [5, 6].

On \mathbb{Z}^d ($d = d_1 + d_2$, $d_1 > 0$, $d_2 \geq 3$), we consider random Hamiltonians of the form

$$(1) \quad H_\omega = -\Delta + V_\omega$$

where

(H0): Let H be a translational invariant Jacobi matrix with exponential off-diagonal decay that is $H = ((h_{\gamma-\gamma'}))_{\gamma, \gamma' \in \mathbb{Z}^d}$ such that,

- $h_{-\gamma} = \overline{h_\gamma}$ for $\gamma \in \mathbb{Z}^d$ and for some $\gamma \neq 0$, $h_\gamma \neq 0$.
- there exists $c > 0$ such that, for $\gamma \in \mathbb{Z}^d$,

$$(2) \quad |h_\gamma| \leq \frac{1}{c} e^{-c|\gamma|}.$$

(H1): V_ω is a random potential concentrated on the sub-lattice $\mathbb{Z}^{d_1} \times \{0\} \subset \mathbb{Z}^d$ of the form

$$(3) \quad V(\gamma_1, \gamma_2) = \begin{cases} \omega_{\gamma_1} & \text{if } \gamma_2 = 0, \\ 0 & \text{if } \gamma_2 \neq 0. \end{cases}, \gamma = (\gamma_1, \gamma_2) \in \mathbb{Z}^{d_1} \times \mathbb{Z}^{d_2} = \mathbb{Z}^d.$$

and $(\omega_{\gamma_1})_{\gamma_1 \in \mathbb{Z}^{d_1}}$ is a family of non trivial i.i.d. bounded random variables.

The operator H_ω is bounded for almost every ω . It is ergodic. So we know there exists Σ the almost sure spectrum of H_ω (see e.g. [4, 8]). Note that the Σ_0 contains the spectrum of H .

Remark 1. *An interesting case which can be brought back to a Hamiltonian of the form (1) with H and V_ω as above is the following.*

Consider Γ , a sub-lattice of \mathbb{Z}^d obtained in the following way $\Gamma = G(\{0\} \times \mathbb{Z}^{d_2})$ where G is a matrix in $GSL_d(\mathbb{Z})$, the d -dimensional special linear group over \mathbb{Z} , i.e. the multiplicative group of invertible matrices with coefficients in \mathbb{Z} and unit determinant. One easily shows that the random operator

$$H_\omega(\Gamma) = -\frac{1}{2}\Delta + \sum_{\gamma \in \Gamma} \omega_\gamma \Pi_\gamma$$

(where Π_γ is the projector onto the vector $\delta_\gamma \in \ell^2(\mathbb{Z}^d)$) is unitarily equivalent to $H + V_\omega$ where V_ω is defined in (3) for h chosen appropriately (see [5]).

For H_ω as in (1) and satisfying (H0) and (H1), one defines the integrated density of surface states (the IDSS in the sequel), say N_s , in the following way (see e.g. [3, 1, 2, 7]): for $\varphi \in \mathcal{C}_0^\infty(\mathbb{R})$, we set

$$(4) \quad (\varphi'', N_s) = \mathbb{E}(\text{tr}(\Pi_1[\varphi(H_\omega) - \varphi(-\Delta)]\Pi_1))$$

where Π_1 is the orthogonal projector on the subspace $\mathbb{C}\delta_0 \otimes \ell^2(\mathbb{Z}^{d_2}) \subset \ell^2(\mathbb{Z}^d)$. Here δ_0 denotes the vector with components $(\delta_{0j})_{j \in \mathbb{Z}^{d_1}}$.

We normalize N_s so that it vanishes below Σ . In [5], we prove that the function N_s is continuous.

We now present our results on the behavior of N_s near the lower edge of Σ (the study near the upper edge is the same). To fix ideas, assume that $0 = \inf \Sigma$.

Definition 2. We say that E , an edge (or boundary) of the spectrum of H_ω , is stable if it is an edge of the spectrum of $H + tV_\omega$ for all $t \in [0, 1]$. If an edge is not stable, we call it a fluctuation edge.

Let ω_- be the infimum of the support of the random variables $(\omega_{\gamma_1})_{\gamma_1}$. Let $h(\theta)$ be the real analytic function

$$h(\theta) = \sum_{\gamma \in \mathbb{Z}^d} h_\gamma e^{i\gamma\theta}.$$

One checks

Proposition 3 ([5]). Write $h(\theta) = h(\theta_1, \theta_2)$ where $\theta = (\theta_1, \theta_2)$, $\theta_1 \in \mathbb{T}^{d_1}$, $\theta_2 \in \mathbb{T}^{d_2}$. Then, 0 is a stable spectral edge if and only if ω_- satisfies condition

$$(5) \quad 1 + \omega_- I_\infty \geq 0 \text{ where } I_\infty := \sup_{\theta_1 \in \mathbb{T}^{d_1}} \int_{\mathbb{T}^{d_2}} \frac{1}{h(\theta_1, \theta_2)} d\theta_2$$

1. THE FLUCTUATION EDGES

We now assume that $\inf \sigma(H) > 0$. In this case, we consider a effective operator \tilde{H} which acts on $\ell^2(\mathbb{Z}^{d_1})$. In Fourier representation this operator is multiplication by the function \tilde{h} given by:

$$(6) \quad \tilde{h}(\theta_1) = \left(\int_{\mathbb{T}^{d_2}} \frac{1}{h(\theta_1, \theta_2)} d\theta_2 \right)^{-1}$$

We either suppose:

(H2): the function $h : \mathbb{T}^d \rightarrow \mathbb{R}$ admits a unique minimum; it is quadratic non-degenerate.

or we assume the weaker hypothesis:

(H2'): the function $\tilde{h} : \mathbb{T}^d \rightarrow \mathbb{R}$ is not constant.

Let P_0 be the common distribution of the random variables $(\omega_{\gamma_1})_{\gamma_1}$ defining the potential (3). We assume:

(H3): P_0 is not trivial and $P_0([\omega_-, \omega_- + \varepsilon]) \geq \varepsilon^k/k$ for some $k > 0$.

We prove

Theorem 4 ([5]). *If (H0) – (H2) and (H3) are satisfied then*

$$\lim_{E \searrow 0} \frac{\ln |\ln(N_s(E))|}{\ln E} = -\frac{d_1}{2}.$$

We have an additional result for low dimension of the surface:

Theorem 5 ([5]). *Assume (H0) – (H2') and (H3) hold. If $d_1 = 1$ then*

$$(7) \quad \lim_{E \searrow 0} \frac{\ln |\ln(N_s(E))|}{\ln E} = -\lim_{E \searrow 0} \frac{\ln(n(E - \omega_-))}{\ln E}$$

where $n(E)$ is the integrated density of states for \tilde{H} .

If $d_2 = 2$, then

$$(8) \quad \lim_{E \searrow 0} \frac{\ln |\ln(N_s(E))|}{\ln(E)} < 0.$$

Both limits (7) and (8) can be computed in terms of the Taylor series of \tilde{h} at its minima (see [5] for details).

2. THE STABLE EDGES WHEN $d_2 \in \{1, 2\}$

In this case, we prove

Theorem 6. *Assume (H0) and (H2) hold. Assume, moreover, that 0 is a stable spectral edge for H_ω . Then,*

- if $d_2 = 1$: $N_s(E) \underset{E \rightarrow 0^+}{\sim} \frac{\text{Vol}(\mathbb{S}^{d_1-1}) \cdot C(h)}{d_1(d_1 + 2)(2\pi)^{d_1}} E^{1+d_1/2}$;
- if $d_2 = 2$: $N_s(E) \underset{E \rightarrow 0^+}{\sim} \frac{2 \text{Vol}(\mathbb{S}^{d_1-1}) \cdot C(h)}{d_1(d_1 + 2)(2\pi)^{d_1}} \frac{E^{1+d_1/2}}{|\log E|}$.

Here, the constant $C(h)$ depends only of the Hessian of h at its minimum.

The striking feature is that to first order these asymptotics are independent of the random potential. The reason for this is that the asymptotics of integrated density of surface states for a constant surface potential near a stable edge does not depend on the value of the potential (to leading order).

3. THE STABLE EDGES WHEN $d_2 \geq 3$

We assume that

(H3): for almost every θ_1 , the function $\theta_2 \mapsto h(\theta_1, \theta_2)$ is not constant.

Consider the embedding $U_2 : \ell^2(\mathbb{Z}^{d_1}) \rightarrow \ell^2(\mathbb{Z}^d)$ defined by $v = U_2(u)$ where

$$(9) \quad v_{\gamma_1, \gamma_2} = u_{\gamma_1} \delta_{\gamma_2, 0} \quad \text{for } u = (u_{\gamma_1})_{\gamma_1 \in \mathbb{Z}^{d_1}}.$$

The embedding U_2 is a partial isometry as $U_2^* U_2 = Id$ on $\ell^2(\mathbb{Z}^{d_1})$. One proves that, under assumptions (H0) – (H3), for almost every ω , the operator $U_2^* H U_2 + V_\omega$ is

positive and the operator $\mathbb{E}((U_2^* H U_2 + V_\omega)^{-1})$ is bounded and positive. Define V_{eff} to be the operator

$$V_{\text{eff}} = \left[\mathbb{E}((U_2^* H U_2 + V_\omega)^{-1}) \right]^{-1} - U_2^* H U_2$$

acting on $\ell^2(\mathbb{Z}^{d_1})$. One proves that the operator V_{eff} acts as a convolution. Let $\theta_1 \mapsto v_{\text{eff}}(\theta_1)$ be the symbol of this operator (i.e. the operator is conjugated to multiplication by this function using the discrete Fourier transform). The function v_{eff} is real analytic on the torus \mathbb{T}^{d_1} . Note that the strict convexity of $x \mapsto 1/x$ for $x > 0$ implies that, for $\theta_1 \in \mathbb{T}^{d_1}$, $\omega_- < v_{\text{eff}}(\theta_1) < \mathbb{E}(\omega_0)$. Our main result is

Theorem 7 ([6]). *Assume that 0 is a stable edge. Under the assumptions (H0) – (H3), one has*

- if $v_{\text{eff}}(0) \neq 0$, then

$$N(E) = \frac{C}{\sqrt{\text{Det} Q}} \frac{v_{\text{eff}}(0)}{1 + v_{\text{eff}}(0) \cdot I} \cdot E^{d/2} (1 + o(1)) \quad \text{as } E \rightarrow 0^+,$$

- if $v_{\text{eff}}(0) = 0$, then

$$N(E) = o(E^{d/2}) \quad \text{as } E \rightarrow 0^+.$$

where

- C is a constant depending only on d_1 and d_2 ;
- Q is the Hessian matrix of h at 0 and $I = \int_{\mathbb{T}^{d_2}} \frac{1}{h(0, \theta_2)} d\theta_2$.

REFERENCES

- [1] Ayham Chahrour. Densité intégrée d'états surfaciques et fonction généralisée de déplacement spectral pour un opérateur de Schrödinger surfacique ergodique. *Helv. Phys. Acta*, 72(2):93–122, 1999.
- [2] Ayham Chahrour and Jaouad Sahbani. On the spectral and scattering theory of the Schrödinger operator with surface potential. *Rev. Math. Phys.*, 12(4):561–573, 2000.
- [3] H. Englisch, W. Kirsch, M. Schröder, and B. Simon. Density of surface states in discrete models. *Phys. Rev. Lett.*, 61(11):1261–1262, 1988.
- [4] W. Kirsch. Random Schrödinger operators. In A. Jensen H. Holden, editor, *Schrödinger Operators*, number 345 in Lecture Notes in Physics, Berlin, 1989. Springer Verlag. Proceedings, Sonderborg, Denmark 1988.
- [5] W. Kirsch and F. Klopp. The band-edge behavior of the density of surfacic states. To appear in MPAG, 2004.
- [6] F. Klopp. The band-edge behavior of the density of surfacic states: an effective model for stable edges. In progress.
- [7] Vadim Kostrykin and Robert Schrader. The density of states and the spectral shift density of random Schrödinger operators. *Rev. Math. Phys.*, 12(6):807–847, 2000.
- [8] L. Pastur and A. Figotin. *Spectra of Random and Almost-Periodic Operators*. Springer Verlag, Berlin, 1992.

Recent results on singular spectrum of Schrödinger operators

ALEXANDER KISELEV

In the recent years, there has been significant interest and progress in studying spectral types of one-dimensional Schrödinger operators with slowly decaying potentials and Stark (constant electric field) operators with rough potentials. Many of the new results concern the operators which can have rich and subtle spectral structure, such as dense imbedded point spectrum, singular continuous spectrum imbedded in the absolutely continuous, or singular continuous spectrum of fixed Hausdorff dimension. New results often involved new technology, such as use of fairly advanced Fourier analysis for studying the asymptotic behavior of solutions or a fruitful interaction of spectral theory methods and methods developed by orthogonal polynomials community. This brief note reviews just a small piece of the big picture consisting of a couple of recent results of the author, partly in collaboration with Michael Christ. The references are far from complete - rather fairly sketchy given the format of the note.

Let us define

$$(1) \quad H_V = -\frac{d^2}{dx^2} + V(x)$$

to be a Schrödinger operator defined on half-axis $\mathbb{R}^+ = (0, \infty)$ with, say, Dirichlet boundary condition at the origin. Let us also denote modified wave operators

$$\Omega_{\pm}^m f = \lim_{t \rightarrow \mp \infty} e^{itH_V} e^{-itH_0 \pm iW(H_0^{1/2}, \mp t)} f$$

for all $f \in L^2(\mathbb{R}^+)$, where existence of the limit has to be established. Here W is given by

$$W(\lambda, t) = -(2\lambda)^{-1} \int_0^{2\lambda t} V(s) dx.$$

Theorem 1. *Assume that the potential $V \in L^p$ with $p < 2$. Then there exist modified wave operators Ω_{\pm}^m . If $\int_0^x V(s) ds$ has a finite limit as x goes to infinity, usual Möller wave operators exist. Moreover, for a.e. k there exist a solution $u(x, k)$ with WKB-type asymptotic behavior as $x \rightarrow \infty$:*

$$(2) \quad u(x, k) = \exp(ikx - \frac{i}{2k} \int_0^x V(s) ds)(1 + o(1)).$$

This theorem appeared in [2]. Classical results on one-dimensional Schrödinger operators with decaying potentials gave L^1 condition. The proof is based on studying the asymptotic behavior of solutions based on almost everywhere convergence results for the multilinear integral operators. The theorem does not hold for $p > 2$ [13, 9]. The theorem is conjectured to be true for $V \in L^2$, but this case is open, and, at least as far as the asymptotic behavior (2) is concerned, presumably very hard. The solution is likely to be related to a nontrivial extension of celebrated Carleson theorem on a.e. convergence of the Fourier series of an L^2 function. It is known, however, that the absolute continuity of the spectrum persists for

$p = 2$. This sharp result is due to Deift and Killip [3], who employ a sum rule to control the spectrum. In the higher dimensions, the slowly decaying perturbations are much less understood. The conjecture of Barry Simon, which is also put forward as one of his fifteen "twenty first century" problems in Schrödinger operators [14], states that the absolutely continuous spectrum is preserved as far as $\int |V(x)|^2(1 + |x|)^{-d+1} dx < \infty$. However, the best general result available is still a classical short range result of Agmon. There are some recent results under mild additional conditions on the oscillation of potential [4, 10], and an interesting result of Bourgain in random case [1].

While in the situation of Theorem 1, the absolutely continuous spectrum fills the whole real axis, the crucial difference with the short range case is that the singular spectrum can also be very rich. Dense imbedded point spectrum is possible due to the results of Naboko and Simon. The set of singular energies (which we define as a Lebesgue measure zero set where the asymptotic behavior (2) fails) can have any Hausdorff dimension ≤ 1 . For $p = 2$, the singular part of the spectral measure can be pretty much arbitrary modulo some normalization conditions, as follows from work of Killip and Simon (see [7] for the discrete case). One of the "twenty first century" problems of Barry Simon has asked whether potentials satisfying $|V(x)| \leq C(1 + |x|)^{-\alpha}$ for $\alpha > 1/2$ can lead to imbedded singular continuous spectrum. Controlling imbedded singular continuous spectrum is difficult, since there is no simple criteria to establish its existence, and the typical approach of proving there is some spectra which cannot be neither pure point nor absolutely continuous [15] does not work. First important progress has been achieved by Denisov [5], who proved that if $V \in L^2$, the singular continuous spectrum may appear (with further beautiful and complete results of Killip and Simon). Our next two theorems provide a sharp answer to the question of a decay rate for which singular continuous spectrum may appear [8].

Theorem 2. *For any function $h(x) \rightarrow \infty$, there exists a potential $V(x)$ such that $|V(x)| \leq \frac{h(x)}{1+x}$ and the singular continuous spectrum of the operator H_V is not empty.*

The theorem also provides, up to the best of my knowledge, the first example where the wave operators exist and cohabit with imbedded singular continuous spectrum, thus leading to the lack of asymptotic completeness. The proof of this theorem is fairly involved; it is based on approximation by operators having imbedded eigenvalues and careful study of the asymptotic behavior of the solutions to establish control over the weights the spectral measure assigns to these eigenvalues. Generalized Prüfer transform and analysis of oscillatory integrals with the nonlinear dependence of phase on the argument function play a key role.

Theorem 3. *If $|V(x)| \leq \frac{C}{1+|x|}$ for some constant C , then the singular continuous spectrum of the operator H_V is empty.*

This theorem shows that the critical threshold is the Coulomb rate of decay and so the construction of the previous theorem is sharp. The proof of the absence of the singular continuous spectrum for potentials decaying at the Coulomb rate

is based on the analysis of approximations where the potential is cut off at a finite scale. The main difficulty lies in the fact that there can be a singular set where the derivative of the spectral measure is infinite and which would be large enough to support the singular continuous spectrum (and can even be dense in $(0, \infty)$!). Therefore, one cannot use some sort of standard resolvent estimates. The technique used involves Gilbert-Pearson subordinacy theory, analysis of the singular set using Fourier transform methods, and a general approximation lemma proved in [6].

REFERENCES

- [1] J. Bourgain, *On random Schrödinger operators on Z^2* , Discrete Contin. Dyn. Syst. **8** (2002), 1–15
- [2] M. Christ and A. Kiselev, *Scattering and wave operators for one-dimensional Schrödinger operators with slowly decaying nonsmooth potentials*, Geom. Funct. Anal. **12** (2002), 1174–1234
- [3] P. Deift and R. Killip, *On the absolutely continuous spectrum of one-dimensional Schrödinger operators with square summable potentials*, Commun. Math. Phys. **203** (1999), 341–347
- [4] S. Denisov, *On the coexistence of absolutely continuous and singular continuous components of the spectral measure for some Sturm-Liouville operators with square summable potentials*, J. Diff Eq. **191** (2003), 90–104
- [5] S. Denisov, *Absolutely continuous spectrum of multidimensional Schrodinger operator*, preprint
- [6] F. Germinet, A. Kiselev and S. Tcheremshantsev, *Transfer matrices and transport for Schrödinger operators*, Annales de l'Institut Fourier, **54** (2004), 787–830
- [7] R. Killip and B. Simon, *Sum rules for Jacobi matrices and their application to spectral theory*, Ann. Math. **158** (2003), 253–321
- [8] A. Kiselev, *Imbedded singular continuous spectrum for Schrödinger operators*, submitted
- [9] S. Kotani and N. Ushiroya, *One-dimensional Schrödinger operators with random decaying potentials*, Commun. Math. Phys. **115** (1988), 247–266
- [10] A. Laptev, S. Naboko and O. Safronov, *Absolutely continuous spectrum of Schrödinger operators with slowly decaying and oscillating potentials*, preprint
- [11] S.N. Naboko, *Dense point spectra of Schrödinger and Dirac operators*, Theor.-math. **68** (1986), 18–28
- [12] T. Kriecherbauer and C. Remling, *Finite gap potentials and WKB asymptotics for one-dimensional Schrödinger operators*, Comm. Math. Phys. **223** (2001), 409–435
- [13] D. Pearson, *Singular continuous measures in scattering theory*, Comm. Math. Phys. **60** (1978), 13–36
- [14] B. Simon, *Some Schrödinger operators with dense point spectrum*, Proc. Amer. Math. Soc. **125** (1997), 203–208
- [15] ———, *Operators with singular continuous spectrum. I. General operators*, Ann. of Math. **141** (1995), 131–145

Dispersive estimates for Schrödinger equations

KENJI YAJIMA

We consider the time decay in L^p spaces of solutions of the initial value problem for three dimensional Schrödinger equations

$$(1) \quad i\partial_t u = (-\Delta + V(x))u, \quad u(0) = \phi \in L^2(\mathbf{R}^3).$$

We assume that the potentials $V(x)$ decay faster than $C\langle x \rangle^{-5/2-\varepsilon}$ at infinity. The operator $H = -\Delta + V$ in the right of (1) is selfadjoint in $L^2(\mathbf{R}^3)$ and the solution of (1) is uniquely given by $u(t) = e^{-itH}\phi$. Let P_c be the orthogonal projection to the continuous spectral subspace for H . Then, $e^{-itH}P_c\phi$ is a scattering solution of (1) and it is now well known that it satisfies the so called L^p - L^q estimates

$$(2) \quad \|e^{-itH}P_c u\|_p \leq C_p t^{-3(\frac{1}{2}-\frac{1}{p})} \|u\|_q, \quad u \in L^2 \cap L^q$$

for $1 \leq q \leq 2 \leq p \leq \infty$, $1/p + 1/q = 1$, provided that 0 is not an eigenvalue nor a resonance of H (Goldberg-Schlag ([5]), see also [7], [1], [15], [15], [16], [13], [10], [12] for earlier and related works). This implies Strichartz inequality and it has been a very useful and important tool for studying linear and nonlinear Schrödinger equations (see e.g. [8]). It is also known that (2) cannot hold for all $2 \leq p \leq \infty$ if H is of exceptional type as it would contradict the local decay estimate of Jensen-Kato[6] or Murata[9].

In this paper, we analyze the behavior as $t \rightarrow \pm\infty$ of scattering solutions of (1) in L^p spaces when 0 is an eigenvalue or/and a resonance of H . We show how (2) is violated and propose a new estimate which replaces (2). To state the main results we introduce some notation. For $1 \leq p, q \leq \infty$, $L^{p,q}$ is the Lorentz space with the norm $\|u\|_{p,q}$. For $\gamma \in \mathbf{R}$, $L^2_\gamma = L^2(\mathbf{R}^3, \langle x \rangle^{2\gamma} dx)$ is the weighted L^2 space. We write $R_0(z) = (H_0 - z)^{-1}$ and $R(z) = (H - z)^{-1}$ for the resolvents of $H_0 = -\Delta$ and H respectively. For $\lambda \in \mathbf{C}$

$$(3) \quad G_0(\lambda)u(x) = \frac{1}{4\pi} \int \frac{e^{i\lambda|x-y|}}{|x-y|} u(y) dy.$$

We have $R_0(\lambda^2) = G_0(\lambda)$ for $\Im\lambda > 0$. The integral kernel of $G_0(\lambda)$ is an entire function of $\lambda \in \mathbf{C}$ and, using its derivatives at $\lambda = 0$, we define

$$(4) \quad D_j u(x) = \frac{1}{4\pi j!} \int |x-y|^{j-1} u(y) dy, \quad j = 0, 1, \dots,$$

so that $G_0(\lambda) = D_0 + i\lambda D_1 + (i\lambda)^2 D_2 + \dots$ at least formally.

For any $1/2 < \gamma < \beta - 1/2$, the operator $D_0 V$ is of Hilbert-Schmidt type in $L^2_{-\gamma}$ and we denote the null space of $1 + D_0 V$ by \mathcal{M} . The space \mathcal{M} is finite dimensional and is independent of $1/2 < \gamma < \beta - 1/2$. All $\phi \in \mathcal{M}$ satisfy the stationary Schrödinger equation $-\Delta\phi(x) + V(x)\phi(x) = 0$ and the converse is also true for $\phi \in L^2_{-\frac{3}{2}}$. The eigenspace \mathcal{E} of H with eigenvalue 0 is therefore a subspace of \mathcal{M} . The function $\phi \in \mathcal{M}$ is in \mathcal{E} if and only if $\langle V, \phi \rangle = 0$ and $\text{codim}_{\mathcal{M}} \mathcal{E} \leq 1$. The sesquilinear form $-(u, Vv)$ is an inner product in \mathcal{M} .

Definition 1. We say H or V is of generic type if $\mathcal{M} = \{0\}$ and is of exceptional type otherwise. H is of exceptional type of the first kind if $\mathcal{M} \neq \{0\}$ and $\mathcal{E} = 0$; of the second kind if $\mathcal{E} = \mathcal{M} \neq \{0\}$; and of the third kind if $\{0\} \subset \mathcal{E} \subset \mathcal{M}$ with strict inclusions. A function $\phi \in \mathcal{M} \setminus \mathcal{E}$ is called a resonance of H .

Any resonance $\phi(x)$ satisfies $\phi(x) - C|x|^{-1} \in L^2$ for a constant $C \neq 0$ and that $\phi \in \mathcal{E}$ may decay as slowly as $C\langle x \rangle^{-2}$. We write P_0 for the orthogonal projection in L^2 onto \mathcal{E} .

When H is of exceptional type of the third kind, we let $\phi_1 \in \mathcal{M}$ be a (uniquely determined) resonance such that $\langle V, \phi_1 \rangle > 0$, $-\langle \phi_1, V\phi_1 \rangle = 1$ and $-\langle \phi_1, V\phi_j \rangle = 0$ for all $\phi_j \in \mathcal{E}$ and define the canonical resonance by $\varphi(x) = \phi_1(x) + P_0VD_2V\phi_1(x)$.

Using $\varphi(x)$, set $a = 4\pi i|\langle V, \varphi \rangle|^{-2}$ and $\zeta(t, x) = e^{i\frac{x^2}{4t}}\varphi(x)$. We define

$$(5) \quad \mu(t, x) = \frac{i}{|x|} \int_0^1 (e^{\frac{i|x|^2}{4t}} - e^{\frac{i\theta^2|x|^2}{4t}})d\theta;$$

$\mu(t)$ is multiplication with $\mu(t, x)$ and $f \otimes g$ is the rank one operator defined by integral kernel $f(x)g(y)$ (not $f(x)g(y)$).

Definition 2. We define the operators $R(t)$ and $S(t)$ respectively by

$$(6) \quad R(t) = \frac{ae^{-i\frac{3\pi}{4}}}{\sqrt{\pi t}}\zeta(t, x) \otimes \zeta(t, x),$$

$$(7) \quad S(t) = \frac{e^{-i\frac{3\pi}{4}}}{\sqrt{\pi t}}(-iP_0VD_3VP_0 + \mu(t)D_2VP_0 + P_0VD_2\mu(t)).$$

When H is of exceptional type of the first or the second kind, we use the same notation, setting, of course, $S(t) = 0$ or $R(t) = 0$ respectively.

We remark that $\zeta(t, x) - \varphi(x)$ and $\mu(t, x)$ are both bounded by

$$(8) \quad C \min\left(\frac{1}{\sqrt{t}}, \frac{1}{|x|}, \frac{|x|}{|t|}\right).$$

As $\phi \in \mathcal{E}$ satisfy $\int V(x)\phi(x)dx = 0$, $(D_2V\phi)(x)$ are bounded and, if $\{\phi_2, \dots, \phi_d\}$ is an orthonormal basis of \mathcal{E} and $w_j(t, x) = \mu(t, x)(D_2V\phi_j)(x)$, $j = 2, \dots, d$, then $w_j(t, x)$ are bounded by (8) and $S(t)$ may be written in the form

$$\frac{e^{\frac{i\pi}{4}}}{\sqrt{\pi t}} \left(\sum_{j,k=2}^d a_{jk}\phi_j \otimes \phi_k + \sum_{j=2}^d (w_j(t) \otimes \phi_j + \phi_j \otimes w_j(t)) \right).$$

Theorem 3. Let V satisfy $|V(x)| \leq C\langle x \rangle^{-\beta}$ for some $\beta > 11/2$. Suppose that H is of exceptional type. Then the following statements are satisfied:

- (i) Estimate (2) holds when $3/2 < q \leq 2 \leq p < 3$ and $1/p + 1/q = 1$.
- (ii) (2) holds when L^3 and $L^{\frac{3}{2}}$ are respectively replaced by $L^{3,\infty}$ and $L^{\frac{3}{2},1}$.
- (iii) When $3 < p \leq \infty$ and $1 \leq q < 3/2$ are such that $1/p + 1/q = 1$, there exists a constant C_{pq} such that for any $u \in L^2 \cap L^q$

$$(9) \quad \|(e^{-itH}P_c - R(t) - S(t))u\|_p \leq C_{pq}t^{-3(\frac{1}{2} - \frac{1}{p})}\|u\|_q.$$

If H is of exceptional type of the first kind, theorem holds under the condition $|V(x)| \leq C\langle x \rangle^{-\beta}$ with $\beta > 9/2$.

Theorem 4. Let V satisfy $|V(x)| \leq C\langle x \rangle^{-\beta}$ for some $\beta > 11/2$. Suppose that H is of exceptional type. Then, for $3 < p \leq \infty$ and $1 \leq q < 3/2$ such that $1/p + 1/q = 1$, there exists a constant C such that

$$(10) \quad \|e^{-itH} P_c u\|_p \leq C t^{-3(\frac{1}{2} - \frac{1}{p})} (\|u\|_q + \|\langle x \rangle^{\frac{6}{q} - 5} u\|_1)$$

for any $u \in L^2 \cap L^q$ which satisfies $\langle \phi, u \rangle = 0$ for all $\phi \in \mathcal{M}$ and $\langle x \rangle^{\frac{6}{q} - 5} u \in L^1$. If H is of exceptional type of the first kind, the same statement holds under a weaker decay condition $|V(x)| \leq C\langle x \rangle^{-\beta}$ with $\beta > 9/2$.

REFERENCES

- [1] G. Artbazar and K. Yajima, *The L^p -continuity of wave operators for one dimensional Schrödinger operators*, J. Math. Sci. Univ. Tokyo **7** (2000), 221–240.
- [2] S. Cuccagna, *Stabilization of solutions to nonlinear Schrödinger equations*, Comm. Pure Appl. Math. **54** (2001), 1110–1145.
- [3] M. B. Erdoğan and W. Schlag, *Dispersive estimates for Schrödinger operators in the presence of a resonance and/or an eigenvalue at zero energy in dimension three I*, preprint (2004).
- [4] M. Goldberg, *Dispersive bounds for the three-dimensional Schrödinger equation with almost critical potentials*, preprint (2004).
- [5] M. Goldberg and W. Schlag, *Dispersive estimates for Schrödinger operators in dimensions one and three*, preprint (2003).
- [6] A. Jensen and T. Kato, *Spectral properties of Schrödinger operators and time-decay of the wave functions*, Duke Math. J. **46** (1979), 583–611.
- [7] J.-L. Journé, A. Soffer and C. D. Sogge, *Decay estimates for Schrödinger operators*, Comm. Pure and Appl. Math **40** (1991), 573–604.
- [8] T. Kato, *On nonlinear Schrödinger equations*, Ann. Inst. H. Poincaré, Phys. Théor. **46** (1987), 113–129.
- [9] M. Murata, *Asymptotic expansions in time for solutions of Schrödinger-type equations*, J. Funct. Anal., **49** (1982), 10–56.
- [10] I. Rodnianski and W. Schlag, *Time decay for solutions of Schrödinger equations with rough and time dependent potentials*, to appear in Invent. Math.
- [11] I. Rodnianski, W. Schlag and A. Soffer, *Dispersive analysis of charge transfer models*, preprint (2003).
- [12] W. Schlag, *Dispersive estimates for Schrödinger operators in dimensions two*, preprint (2004).
- [13] R. Weder, *L^p - $L^{p'}$ estimates for the Schrödinger equations on the line and inverse scattering for the nonlinear Schrödinger equation with a potential*, J. Funct. Anal. **170** (2000), 37–68.
- [14] K. Yajima, *The $W^{k,p}$ -continuity of wave operators for Schrödinger operators*, J. Math. Soc. Japan **47** (1995), 551–581.
- [15] K. Yajima, *The $W^{k,p}$ -continuity of wave operators for Schrödinger operators III*, J. Math. Sci. Univ. Tokyo **2** (1995), 311–346.
- [16] K. Yajima, *L^p -boundedness of wave operators for two dimensional Schrödinger operators*, Commun. Math. Phys. **208** (1999), 125–152.

Classical and Quantum Mechanics for a Particle in a Long-Range Magnetic Field

IRA HERBST

This talk is about some work of Horia Cornean, Erik Skibsted, and I in progress [CHS2] concerning the dynamics of a charged particle moving in a plane subject to a magnetic field which is homogeneous of degree -1 . The work the talk is drawn from also deals with electric forces with the same homogeneity, but for simplicity, here we set the electric potential equal to zero. We analyze the classical and quantum dynamics of this system for large time with the objective to prove asymptotic completeness in quantum mechanics with some simple appropriate approximate dynamics.

Thus consider a magnetic field of the form

$$B = \frac{b(\theta)}{r},$$

where (r, θ) are the polar coordinates of a point in R^2 . We always assume that b is smooth (and periodic of period 2π). Introducing the velocities

$$\begin{aligned}\rho &= \frac{dr}{dt}, \\ \eta &= \frac{rd\theta}{dt},\end{aligned}$$

and the new time τ given by

$$\frac{d\tau}{dt} = \frac{1}{r},$$

we can write the equations of motion of the particle in a reduced phase space as

$$\begin{aligned}\frac{d\rho}{d\tau} &= \eta(\eta + b(\theta)), \\ \frac{d\eta}{d\tau} &= -\rho(\eta + b(\theta)), \\ \frac{d\theta}{d\tau} &= \eta.\end{aligned}$$

Introducing the angle ϕ by

$$\begin{aligned}\rho &= \sqrt{2E} \sin \phi, \\ \eta &= \sqrt{2E} \cos \phi,\end{aligned}$$

where E is the conserved kinetic energy, the first two equations above become

$$\frac{d\phi}{d\tau} = \sqrt{2E} \cos \phi + b(\theta),$$

which shows that the reduced classical phase space at energy E is a 2-torus. Note that r can be found once ρ is known:

$$r = r_0 e^{\int_0^\tau \rho(\tau') d\tau'}.$$

The case of $b < 0$ was treated in [CHS], where it was shown that in classical mechanics, above a certain energy E_d there is an attracting periodic orbit on the torus which attracts all orbits except for another periodic orbit which only lives for a finite amount of real time. Asymptotic completeness was proved in quantum mechanics above E_d using a semiclassical approximate dynamics based on this attracting periodic orbit. Below E_d nothing is known. But in the case where b is a non-zero constant the Hamiltonian has dense point spectrum [CFKS] below E_d .

We consider below mostly the classical mechanics of the model and only mention any difficulties that arise in quantum mechanics. One of these difficulties arises immediately when we consider the classical observable

$$A_1 = \rho - \int_0^\theta b(\theta') d\theta'.$$

Note that

$$A_1(\tau_2) - A_1(\tau_1) = \int_{\theta(\tau_1)}^{\theta(\tau_2)} \eta(\tau)^2 d\tau.$$

The problem with this observable is that unless the “flux” $\int_0^{2\pi} b(\theta) d\theta = 0$, it is not a function on the torus (but rather on a covering space of the torus), so it does not have a good quantization. Let

$$A_2 = -\rho\eta b(\theta).$$

Then for bounded E , we have that for C large enough

$$\frac{d(CA_1 + A_2)}{d\tau} \geq Eb^2 + \eta^2.$$

Let us assume that the flux is ≤ 0 , and that b has zeros but all are non-degenerate. Then with some additional work, it follows that either

(1)

$$\theta(\tau) \rightarrow \infty,$$

or

(2)

$$\lim_{\tau \rightarrow \infty} [\theta(\tau) - \theta_0]^2 + \eta(\tau)^2 = 0,$$

where $b(\theta_0) = 0$.

Consider the fixed point in (2) where at $\tau = \infty$, $\rho = \sqrt{2E}$. Then this classical channel has a corresponding quantum channel if and only if the fixed point is a sink on the torus [HS2] (this corresponds to $b'(\theta_0) > 0$). Otherwise the fixed point has a corresponding stable manifold, but there are no states for which θ approaches θ_0 in quantum mechanics. The fixed points in (2) for which $\rho = -\sqrt{2E}$ correspond to orbits which hit the origin in finite time and have no analog in quantum mechanics. The $\rho = +\sqrt{2E}$ quantum channels may be described by an approximate dynamics

as in [HS1] and asymptotic completeness proved, at least if the energy is high enough.

If (1) obtains, then aside from a finite set of energies, if the orbit does not collapse at the origin, it is attracted to a periodic orbit in the reduced phase space. There is also a corresponding channel in quantum mechanics. Asymptotic completeness can be proved there using a simple semiclassical approximate dynamics as in [CHS].

The analysis in [CHS2] consists first of a detailed description of the classical dynamics of this model. This alone is very non-trivial. Where quantum observables with positive Heisenberg derivatives are available to prove appropriate smoothness estimates, they are used. But there seem to be situations where not enough of these observables are available, and we supplement the analysis with “propagation of decay” estimates as in the usual propagation of singularities theorems. This idea (with x and ξ reversed) of using propagation of singularity theorems was introduced into scattering theory by Melrose in [M] and used for example in [HMV] in work closely related to [HS1].

REFERENCES

- [CFKS] Cycon, H., Froese, R., Kirsch, W., Simon, B.: *Schrödinger Operators with Applications to Quantum Mechanics and Global Geometry*, Berlin: Springer Verlag, 1987.
- [CHS] Cornean, H., Herbst, I., Skibsted, E.: Spiraling attractors and quantum dynamics for a class of long-range magnetic fields, MaPhySto preprint No. 36, 2002.
- [CHS2] _____: Classical and quantum dynamics for 2D-electromagnetic potentials asymptotically homogeneous of degree zero, in progress.
- [HS1] Herbst, I., Skibsted, E.: Quantum scattering for potentials independent of $|x|$: Asymptotic completeness for high and low energies, *Comm. P.D.E.*, **24**, 547–610, 2004.
- [HS2] _____: Absence of quantum states corresponding to unstable classical channels: Homogeneous potentials of degree zero, MaPhySto preprint no. 23, 2003.
- [HMV] Hassell, A., Melrose, R., Vasy, A.: Spectral and scattering theory for symbolic potentials of order zero, *Advances in Mathematics* **181**, 1–87, 2004.
- [M] Melrose, R.: Spectral and scattering theory for the Laplacian on asymptotically Euclidean spaces, *Spectral and Scattering Theory* (Sanda, 1992) (M. Ikawa, ed.), Marcel Dekker (1994), 85–130.

Some variational principles for relativistic energy functionals

JEAN-MARIE BARBAROUX

(joint work with V. Bach, M. Esteban, W. Farkas, B. Helffer, E. Séré and H. Siedentop)

We give here some connections between two models describing the energy of a system of relativistic particles in the field of a pointwise fixed nucleus: The Dirac-Fock equations [4, 7], derived from the so-called Dirac-Fock functional \mathcal{E}^{DF} , and the electron/positron field functional \mathcal{E}^{e^-/e^+} (see (2) below) derived from a simple no photon QED formal Hamiltonian, in the generalized Hartree-Fock approximation [1, 3].

The Hamiltonien for one electron in the field of a nucleus of charge eZ is given by the Coulomb-Dirac operator

$$(1) \quad D_Z := \boldsymbol{\alpha} \cdot \frac{1}{i} \nabla + m\beta - e^2 \frac{Z}{|\mathbf{x}|} \quad \text{on } \mathfrak{H} := L^2(\mathbb{R}^3) \otimes \mathbb{C}^4,$$

where e^2 is the Sommerfeld fine structure constant, and α and β are the 4×4 Dirac matrices. Here, we assume $e^2 Z \in [0, \sqrt{3}/2)$. For $Z = 0$, D_0 is the free Dirac operator. In the following we will also need the Coulomb-Dirac operator written in another system of units:

$$D_c := c\boldsymbol{\alpha} \cdot \frac{1}{i} \nabla + mc^2\beta - \frac{1}{|\mathbf{x}|},$$

where c is the speed of light. Let \mathfrak{H}_+ be a closed subspace of \mathfrak{H} , and define Λ_+ to be the orthogonal projection onto \mathfrak{H}_+ , $\Lambda_- := 1 - \Lambda_+$ and $\mathfrak{H}_- := \Lambda_- \mathfrak{H} = (\mathfrak{H}_+)^{\perp}$. We construct the following variational sets

$$\begin{aligned} S(\mathfrak{H}_+) &= \{ \gamma \in \mathfrak{S}_1(\mathfrak{H}) \mid \gamma = \gamma^*, \operatorname{tr}(|D_0|^{\frac{1}{2}} |\gamma| |D_0|^{\frac{1}{2}}) < \infty, -\Lambda_- \leq \gamma \leq \Lambda_+, \}, \\ S_N(\mathfrak{H}_+) &= \{ \gamma \in S(\mathfrak{H}_+) \mid \operatorname{tr} \gamma = N \}, \\ T_N(\mathfrak{H}_+) &= \{ \gamma \in S(\mathfrak{H}_+) \mid \operatorname{tr} \gamma = N, \Lambda_- \gamma \Lambda_+ = 0 \}, \end{aligned}$$

where $\mathfrak{S}_1(\mathfrak{H})$ denotes the space of trace class operators on \mathfrak{H} . For $\gamma \in S(\mathfrak{H}_+)$, $\operatorname{tr} \gamma$ is the charge of the system (corresponding to an electronic charge $-e \operatorname{tr} \gamma$). The electron/positron field functional we consider is

$$(2) \quad \begin{aligned} \mathcal{E}^{e^-/e^+} : S(\mathfrak{H}_+) &\rightarrow \mathbb{R} \\ \gamma &\mapsto \operatorname{tr}(D_Z \gamma) + \frac{e^2}{2} \int \frac{\overline{\rho_\gamma(\mathbf{x})} \rho_\gamma(\mathbf{y})}{|\mathbf{x}-\mathbf{y}|} d\mathbf{x} d\mathbf{y} - \frac{e^2}{2} \int \frac{\overline{\gamma(x,y)} \gamma(x,y)}{|\mathbf{x}-\mathbf{y}|} dx dy \end{aligned}$$

where $x = (\mathbf{x}, \sigma)$ and $y = (\mathbf{y}, \tau)$ are in $\mathbb{R}^3 \times \{1, 2, 3, 4\}$, $\gamma(x, y)$ is the kernel of γ and $\rho_\gamma(\mathbf{x}) = \sum_{\sigma=1}^4 \gamma((\mathbf{x}, \sigma), (\mathbf{x}, \sigma))$.

Our first result states that without any constraints on the charge, the most stable projection Λ_+ , i.e., the one yielding the highest ground state energy, is given by the projection onto the positive spectral subspace of the Coulomb-Dirac operator.

We denote by \mathfrak{T} the set of all closed subspace \mathfrak{H}_+ of \mathfrak{H} such that the orthogonal projections Λ_{\pm} onto \mathfrak{H}_{\pm} leave $\mathcal{D}(D_Z)$ invariant.

Theorem 1. [1] *Consider D_Z with values of $e, Z \geq 0$ such that $e^2 \leq 4(1-2e^2 Z)/\pi$. We have*

$$(3) \quad \sup_{\mathfrak{H}_+ \in \mathfrak{T}} \inf_{\gamma \in S(\mathfrak{H}_+)} \mathcal{E}^{e^-/e^+}(\gamma) = \inf_{\gamma \in S(\chi_{(0,+\infty)}(D_Z))} \mathcal{E}^{e^-/e^+}(\gamma) = \mathcal{E}^{e^-/e^+}(0) = 0.$$

Moreover, the supremum in (3) is attained only for $\mathfrak{H}_+ = \chi_{(0,+\infty)}(D_Z)$.

We now discuss the case of systems with fixed total charge $N \in \mathbb{N}$. For that purpose, we first need to define Dirac-Fock operators. For $\delta \in F := \{ \delta \in \mathfrak{S}_1(\mathfrak{H}) \mid \delta =$

δ^* , $\text{tr}(|D_0|^{1/2} |\delta| |D_0|^{1/2}) < \infty$, we construct the associated Dirac-Fock operator $D^{(\delta)}$ as

$$\begin{aligned} D^{(\delta)}\psi(x) &= D_Z\psi(x) + e^2W^{(\delta)}\psi(x) \\ &= D_Z\psi(x) + e^2 \int \frac{\rho_\delta(\mathbf{y})}{|\mathbf{x} - \mathbf{y}|} d\mathbf{y} \psi(x) + e^2 \int \frac{\delta(x, y)\psi(y)}{|\mathbf{x} - \mathbf{y}|} dy, \end{aligned}$$

where $\delta(x, y)$ is the kernel of δ and $\rho_\delta(\mathbf{x}) = \sum_{\sigma=1}^4 \delta((\mathbf{x}, \sigma); (\mathbf{x}, \sigma))$. Here, $W^{(\delta)}$ is the mean field Dirac-Fock potential created by the N electrons in the state δ .

We define the associated one-electron space

$$\mathfrak{H}_+^{(\delta)} = \Lambda_+^{(\delta)} \mathfrak{H},$$

where

$$\Lambda_+^{(\delta)} = \chi_{(0, +\infty)}(D^{(\delta)}).$$

As argued in [6], the equality (3) suggests to explore a max-min variational problem similar to (3) in the case of atomic systems with prescribed electronic charge $e(Z - N)$, in order to find the ground state energy:

$$(4) \quad \sup_{\delta \in F} \inf_{\gamma \in T_N(\mathfrak{H}_+^{(\delta)})} \mathcal{E}^{e^-/e^+}(\gamma).$$

The next result shows the existence of solutions for the minimization procedure in (4), and gives the properties of the minimizers.

Theorem 2. [3] *Let $0 \leq \delta \in F$ and assume $Z \geq N \in \mathbb{N}$ such that*

$$e^2\pi(N + 1/4)/(1 - 2e^2Z - 4e^2N) < 1.$$

Then $\mathcal{E}^{e^-/e^+}|_{T_N(\mathfrak{H}_+^{(\delta)})}$ has a minimizer in $T_N(\mathfrak{H}_+^{(\delta)})$, and each minimizer γ^0 is equal to the spectral projection onto the N first eigenvalues of the projected Dirac-Fock operator $\Lambda_+^{(\delta)} D^{(\gamma^0)} \Lambda_+^{(\delta)}$: there exist $\varphi_1^0, \varphi_2^0, \dots, \varphi_N^0$ in $\Lambda_+^{(\delta)} \mathfrak{H} \cap (H^{1/2}(\mathbb{R}^3) \otimes \mathbb{C}^4)$, normalized and orthogonal, and $(\epsilon_i^0)_{i=1, \dots, N}$ in $(0, m)$ such that

$$\gamma^0 = \sum_{i=1}^N |\varphi_i^0\rangle \langle \varphi_i^0|,$$

and

$$\Lambda_+^{(\delta)} D^{(\gamma^0)} \Lambda_+^{(\delta)} \varphi_i^0 = \epsilon_i^0 \varphi_i^0, \quad i = 1, \dots, N,$$

where $(\epsilon_i^0)_{i=1, \dots, N}$ are the N lowest eigenvalues in $(0, m)$ of $\Lambda_+^{(\delta)} D^{(\gamma^0)} \Lambda_+^{(\delta)}$.

We can now compare the max-min procedure (4) with the solutions obtained in [4, 5, 7] by solving the Dirac-Fock equations. We discuss here the nonrelativistic limit case, i.e., with D_Z replaced by D_c , $c \ll 1$, and for $\delta \in F$, $D^{(\delta)} = D_c + e^2W^{(\delta)}$.

Let $\lambda_1 < \lambda_2 < \dots$ be the ordered (positive) eigenvalues of the Coulomb-Dirac operator D_c and let $N_i := \dim(\text{Ker}(D_c - \lambda_i))$ be the dimension of the associated eigenspaces.

Theorem 3. [2][Close to the linear closed shells case]

Let N be the number of electrons, and assume $c \gg 1$ and $e^2 \ll 1$. If $N = \sum_{i=1}^K N_i$ (closed shells), then the variational problem (4) is attained by the self-consistent pair (γ^0, γ^0) , where $\gamma^0 = \sum_{i=1}^N |\varphi_i^0\rangle\langle\varphi_i^0|$, with

$$\Lambda_+^{(\gamma^0)} = \chi_{(0,+\infty)}(D^{(\gamma^0)}),$$

and

$$\left(D_c + e^2 W^{(\gamma^0)}\right) \varphi_i^0 = \epsilon_i^0 \varphi_i^0, \quad \epsilon_i^0 \in (0, m), \quad i = 1, \dots, N,$$

i.e., the N -uple $(\varphi_1^0, \dots, \varphi_N^0)$ is solution of the self-consistent Dirac-Fock equations. Moreover, it is the ground state solution of the Dirac-Fock equations in the sense that it yields the smallest Dirac-Fock energy among the solutions of the Dirac-Fock equations: for any solution (ψ_1, \dots, ψ_N) of the self-consistent Dirac-Fock equations, we have

$$\begin{aligned} \sum_{i=1}^N (\psi_i, D_c \psi_i) + \frac{e^2}{2} \sum_{i \neq j} \left(\int \frac{|\psi_i(x)|^2 |\psi_j(y)|^2}{|\mathbf{x} - \mathbf{y}|} dx dy - \int \frac{\psi_i(x) \overline{\psi_j(y)} \psi_i(y) \psi_j(x)}{|\mathbf{x} - \mathbf{y}|} dx dy \right) \\ \leq \sum_{i=1}^N (\varphi_i^0, D_c \varphi_i^0) + \frac{e^2}{2} \sum_{i \neq j} \left(\int \frac{|\varphi_i^0(x)|^2 |\varphi_j^0(y)|^2}{|\mathbf{x} - \mathbf{y}|} dx dy \right. \\ \left. - \int \frac{\varphi_i^0(x) \overline{\varphi_j^0(y)} \varphi_i^0(y) \varphi_j^0(x)}{|\mathbf{x} - \mathbf{y}|} dx dy \right) = \mathcal{E}^{e^-/e^+}(\gamma^0). \end{aligned}$$

REFERENCES

- [1] V. Bach, J.-M. Barbaroux, B. Helffer, and H. Siedentop. *On the stability of the relativistic electron-positron field*, Comm. Math. Phys. **201** (1999), 445–460.
- [2] J.-M. Barbaroux, M. J. Esteban, and E. Séré. *Some connections between Dirac-Fock and electron-positron Hartree-Fock*, to appear in Ann. Henri Poincaré.
- [3] J.-M. Barbaroux, W. Farkas, B. Helffer, and H. Siedentop. *On the Hartree-Fock equations of the electron-positron field*, to appear in Comm. Math. Phys.
- [4] Maria J. Esteban and Eric Séré. *Solutions of the Dirac-Fock equations for atoms and molecules*, Comm. Math. Phys. **203**(3) (1999), 499–530.
- [5] M. J. Esteban and E. Séré. *Nonrelativistic limit of the Dirac-Fock equations*, Ann. Henri Poincaré **2**(5) (2001), 941–961.
- [6] Marvin H. Mittleman. *Theory of relativistic effects on atoms: Configuration-space Hamiltonian*, Phys. Rev. A **24**(3) (1981), 1167–1175.
- [7] Eric Paturel. *Solutions of the Dirac-Fock equations without projector*, Ann. Henri Poincaré **1**(6)(2000), 1123–1157.

On the mathematical model of the irreversible quantum graph

MIKHAIL SOLOMYAK

Some time ago the physicist Uzy Smilansky suggested a mathematical model which he called “Irreversible quantum graph”. In this model an interaction between the Laplacian on a metric graph Γ and the harmonic oscillator in an “outer

space” is studied. The interaction is introduced by means of the boundary condition of a specific type. This condition involves the coupling parameter $\alpha \geq 0$ which expresses the strength of interaction. For $\alpha = 0$ the interaction is absent.

In the mathematical language the problem consists in the study of the spectral properties of a self-adjoint operator \mathbf{A}_α in the Hilbert space $L^2(\Gamma \otimes \mathbb{R})$. For simplicity, we consider the case when $\Gamma = \Gamma_d$, i.e. the star graph with d edges, each of infinite length, emanating from the only vertex o , the root of the tree. The operator is defined by the differential expression

$$\mathcal{A}U(x, q) = -U''_{xx} + \frac{1}{2}(-U''_{qq} + q^2U), \quad x \in \Gamma \setminus \{o\}, \quad q \in \mathbb{R},$$

and the condition

$$[U'_x](o, q) = \alpha q U(o, q), \quad \forall q \in \mathbb{R}.$$

Here $[U'_x]$ stands for the combination of derivatives appearing in the classical Kirchhoff condition.

On the first glance, this can be reduced to a typical problem of Perturbation Theory for operators defined via their quadratic forms. However, the perturbation turns out to be too strong: it is only bounded but not compact with respect to the unperturbed quadratic form. For this reason, the standard approaches do not apply, and the character of results is rather unusual. Their most important feature is a “phase transition” at the value $\alpha = d/\sqrt{2}$ of the parameter: the spectral properties of the operator A_α for $\alpha\sqrt{2} < d$ and for $\alpha\sqrt{2} > d$ are quite different.

For $\alpha = 0$ separation of variables shows that

$$\sigma(\mathbf{A}_0) = \sigma_{a.c.}(\mathbf{A}_0) = [1/2, \infty);$$

$$\mathbf{m}_{a.c.}(\lambda; \mathbf{A}_0) = dn \quad \text{for } \lambda \in (n - 1/2, n + 1/2), \quad n \in \mathbb{N}.$$

Here $\mathbf{m}_{a.c.}(\lambda; \cdot)$ stands for the multiplicity function for a self-adjoint operator.

The following results describe the picture for $\alpha > 0$.

1. Let $0 < \alpha\sqrt{2} < d$. Then the operator \mathbf{A}_α is positive definite;

$$\sigma_{a.c.}(\mathbf{A}_\alpha) = \sigma_{a.c.}(\mathbf{A}_0) = [1/2, \infty),$$

and the similar equality is satisfied for the multiplicity function.

The operator has no eigenvalues $\geq 1/2$. The spectrum on $(0, 1/2)$ is non-empty and finite. The number $N_-(1/2; \mathbf{A}_\alpha)$ of eigenvalues satisfies the asymptotic formula

$$N_-(1/2; \mathbf{A}_\alpha) \sim \frac{1}{4\sqrt{2}(\mu(\alpha) - 1)}, \quad \mu(\alpha) = \frac{d}{\alpha\sqrt{2}} \quad \text{as } \alpha\sqrt{2} \nearrow d$$

2. Let $\alpha\sqrt{2} \geq d$. Then the operator has no eigenvalues;

$$\sigma_{a.c.}(\mathbf{A}_\alpha) = \begin{cases} [0, \infty), & \alpha\sqrt{2} = d; \\ \mathbb{R}, & \alpha\sqrt{2} > d. \end{cases}$$

$$\mathbf{m}_{a.c.}(\lambda; \mathbf{A}_\alpha) = 1 + \mathbf{m}_{a.c.}(\lambda; \mathbf{A}_0).$$

The results for $\alpha\sqrt{2} \geq d$ were obtained in cooperation with S.N. Naboko.

So, at $\alpha\sqrt{2} = d$ the point spectrum disappears and a new branch of $\sigma_{a.c.}$ arises.

For the proof we use the variational techniques (case $\alpha\sqrt{2} < d$) and the techniques of operator-valued analytic functions (case $\alpha\sqrt{2} \geq d$). In both cases Jacobi matrices arise and play the decisive role in the analysis.

REFERENCES

- [1] *Irreversible quantum graphs*, *Waves in Random Media*, **14** (2004), 143 – 153.
- [2] M.Solomyak, *On a differential operator appearing in the theory of irreversible quantum graphs*, *Waves in Random Media* **14** (2004), 173 – 185.
- [3] M.Solomyak *On the discrete spectrum of a family of differential operators*, *Func. Anal. and Appl.*, **38** (2004), 3, 217 – 223.
- [4] S.N. Naboko and M.Solomyak *On the absolutely continuous spectrum of a family of operators appearing in the theory of irreversible quantum systems*, *Proc. LMS* (submitted).

A Lieb–Thirring Inequality and an Isoperimetric Problem for Closed Curves in \mathbb{R}^2

RAFAEL D. BENGURIA

(joint work with Michael Loss)

The Lieb–Thirring inequalities [8] play a crucial role in the proof of the stability of matter [9]. Let $H = -\Delta + V$ be the Schrödinger operator acting on $L^2(\mathbb{R}^n)$, $n \geq 1$ and denote by $e_1 \leq e_2 \leq \dots < 0$ the negative eigenvalues of H . The Lieb–Thirring inequalities are given by

$$(1) \quad \sum_{j \geq 1} |e_j|^\gamma \leq L_{\gamma,n} \int_{\mathbb{R}^n} V_-(x)^{\gamma+n/2} dx,$$

where $V_-(x) \equiv \max(-V(x), 0)$ is the negative part of the potential. The above inequalities hold for $\gamma \geq 1/2$ when $n = 1$, for $\gamma > 0$ when $n = 2$, and for $\gamma \geq 0$ for $n \geq 3$. The sharp constants for the Lieb–Thirring are known for any $n \geq 1$ when $\gamma \geq 3/2$ and also in the case $n = 1$, $\gamma = 1/2$. See e.g., [6] and references therein for the best constants to date. The sharp constants for the one dimensional Lieb–Thirring inequalities with exponent $\gamma \in (1/2, 3/2)$ are still not known. Lieb and Thirring have conjectured [10] that the sharp constants for this range of exponents should be attained by potentials having only one bound state, and therefore,

$$(2) \quad L_{\gamma,1} \equiv L_{\gamma,1}^1 = \frac{1}{\sqrt{\pi}} \frac{1}{\gamma - 1/2} \frac{\Gamma(\gamma + 1)}{\Gamma(\gamma + 1/2)} \left(\frac{\gamma - 1/2}{\gamma + 1/2} \right)^{\gamma+1/2}$$

([7, 10]).

We have recently shown [1] that there is a connection between this conjecture for $\gamma = 1$ and $n = 1$ and an (still open) isoperimetric inequality for smooth, closed curves, with positive curvature in \mathbb{R}^2 .

Let's denote by C a smooth closed curve in the plane, of length 2π , with positive curvature $\kappa(s)$, and let

$$(3) \quad H(C) \equiv -\frac{d^2}{ds^2} + \kappa^2$$

acting on $L^2(C)$ with periodic boundary conditions. Here s denotes arc-length. Let $\lambda_1(C)$ the lowest eigenvalue of $H(C)$. It has been conjectured (see e.g., [3, 4]), that

$$(4) \quad \lambda_1(C) \geq 1,$$

with equality for a one parameter family of curves that includes the circle.

In recent years several authors have obtained isoperimetric inequalities for the lowest eigenvalues of a variant of $H(C)$. Consider the Schrödinger operator

$$(5) \quad H_g(C) \equiv -\frac{d^2}{ds^2} + g\kappa^2$$

defined on $L^2(C)$ with periodic boundary conditions. As before, C denotes a closed curve in \mathbb{R}^2 with positive curvature κ , and length 2π . If $g < 0$, the lowest eigenvalue of $H_g(C)$, say $\lambda_1(g, C)$ is uniquely maximized when C is a circle [2]. When $g = -1$, the second eigenvalue, $\lambda_2(-1, C)$ is uniquely maximized when C is a circle [5]. If $0 < g \leq 1/4$, $\lambda_1(g, C)$ is uniquely minimized when C is a circle [3]. It is an open problem to determine the curve C that minimizes $\lambda_1(g, C)$ in the cases, $1/4 < g \leq 1$, and $g < 0, g \neq -1$. If $g > 1$ the circle is not a minimizer for $\lambda_1(g, C)$ (see, e.g., [3, 4] for more details on the subject).

Our main result [1] is the following theorem:

Theorem 1. *Suppose that the Schrödinger operator $H = -d^2/dx^2 + V$, acting on $L^2(\mathbb{R})$, has only two negative eigenvalues, say $e_1 < e_2 < 0$. Then, if the isoperimetric inequality (4) holds, we have*

$$(6) \quad |e_1| + |e_2| \leq L_{1,1}^1 \int_{\mathbb{R}} V_-(x)^{3/2} dx.$$

ACKNOWLEDGEMENT

We thank Alexander Sobolev and Timo Weidl for their kind invitation to participate in the workshop *Spectral Analysis of Partial Differential Equations* at the Mathematisches Forschungsinstitut at Oberwolfach. This work has been supported by FONDECYT (Chile) projects 102-0844 and 702-0844.

REFERENCES

- [1] R. D. Benguria, and M. Loss, *Connection between the Lieb–Thirring conjecture for Schrödinger operators and an isoperimetric problem for ovals on the plane*, Contemporary Mathematics **362**(2004), 53–61.
- [2] P. Duclos, and P. Exner, *Curvature-induced bound states in quantum waveguides in two and three dimensions*. Rev. Math. Phys. **7** (1995) 73–102.
- [3] P. Exner, E.M. Harrell, and M. Loss, *Optimal eigenvalues for some Laplacians and Schrödinger operators depending on curvature*. Mathematical results in quantum mechanics (Prague 1998), Oper. Theory Adv. Appl. **108** (1999), 47–58.

- [4] E.M. Harrell, *Gap estimates for Schrödinger operators depending on curvature*, talk delivered at the 2002 UAB International Conference on Differential Equations and Mathematical Physics. Available electronically at <http://www.math.gatech.edu/harrell/>
- [5] E.M. Harrell, and M. Loss, *On the Laplace operator penalized by mean curvature*. Commun. Math. Phys., **195** (1998) 643–650.
- [6] D. Hundertmark, A. Laptev, and T. Weidl, *New bounds on the Lieb-Thirring constants*. Invent. Math. **140** (2000) 693–704.
- [7] Joseph B. Keller, *Lower bounds and isoperimetric inequalities for eigenvalues of the Schrödinger equation*. J. Mathematical Phys. **2** (1961), 262–266.
- [8] E. H. Lieb, *Lieb–Thirring Inequalities* in Encyclopaedia of Mathematics, Suppl. II, Kluwer, Dordrecht 2000, pp. 311–312.
- [9] E.H. Lieb, and W. Thirring, *Bounds for the kinetic energy of fermions which proves the stability of matter*. Phys. Rev. Lett. **35** (1975) 687. Errata: PRL 35, (75) 1116.
- [10] E.H. Lieb, and W. Thirring, *Inequalities for the moments of the eigenvalues of the Schrödinger hamiltonian and their relation to Sobolev inequalities*, in Studies in Mathematical Physics, Essays in Honor of Valentine Bargmann, edited by E.H. Lieb, B. Simon and A.S. Wightman, Princeton University Press, Princeton, NJ 1986, pp. 269–303.

Accurate estimates for magnetic bottles in connection with superconductivity

B. HELFFER

(joint work with S. Fournais)

In this talk, which refers to [FoHe2], we consider a magnetic Schrödinger operator with Neumann boundary conditions in a smooth, bounded domain Ω . We are interested in finding an accurate description of the eigenvalues near the bottom of the spectrum. In particular, we will improve estimates given in [HeMo] in the case of constant magnetic field.

Apart from its intrinsic mathematical interest, this question is important for applications to superconductivity. Precise knowledge of the lowest eigenvalues of this magnetic Schrödinger operator is crucial for a detailed description of the nucleation of superconductivity (on the boundary) for superconductors of Type II and for accurate estimates of the critical field H_{C_3} . We refer the reader to the works of Bernoff-Sternberg [BeSt] who are the first to propose the main conjecture on the basis of formal constructions of quasimodes, Lu-Pan [LuPa1, LuPa2, LuPa3] and Del Pino-Felmer-Sternberg [PiFeSt] for further discussion of this subject.

The domain $\Omega \subset \mathbb{R}^2$ is supposed to be smooth, bounded and simply connected. Points (x_1, x_2) in \mathbb{R}^2 are denoted by x . At each point x of the boundary, we denote by $\nu(x)$ the interior unit normal vector to the boundary of Ω . We define the magnetic Neumann operator \mathcal{H} by

$$(1) \quad \mathcal{D}(\mathcal{H}) \ni u \mapsto \mathcal{H}u = \mathcal{H}_{h,\Omega}u = (-ih\nabla_x - A(x))^2u.$$

Here $A(x) = (-x_2/2, x_1/2)$, so that $\text{curl } A = 1$, and the domain $\mathcal{D}(\mathcal{H})$ of the operator \mathcal{H} is defined by

$$\mathcal{D}(\mathcal{H}) = \{u \in H^2(\Omega) \mid \nu \cdot (-ih\nabla_x - A(x))u|_{\partial\Omega} = 0\}.$$

The case of the half-plane, $\Omega = \mathbb{R} \times \mathbb{R}_+$, will be important for fixing notations and determining the main term of the asymptotics. After a gauge transformation and a partial Fourier transformation, we get, when $h = 1$, the family of models on the half-line:

$$(2) \quad H^{N,\xi} = D_t^2 + (t + \xi)^2 ,$$

on $L^2(\mathbb{R}_+)$ and with Neumann boundary conditions at $t = 0$. Let $\hat{\mu}^{(1)}(\xi)$ be the lowest eigenvalue of $H^{N,\xi}$. Then $\xi \mapsto \hat{\mu}^{(1)}(\xi)$ has a unique minimum Θ_0 attained at a point that we will denote by ξ_0 . The corresponding unique positive, normalized eigenfunction of H^{N,ξ_0} will be denoted by u_0 . We also introduce :

$$(3) \quad C_1 = \frac{u_0^2(0)}{3} .$$

The main result in [FoHe2] gives the asymptotic expansion of the lowest eigenvalues of \mathcal{H} .

Theorem 1. *Suppose that Ω is a smooth bounded domain, that its curvature $\partial\Omega \ni s \mapsto \kappa(s)$ at the boundary has a unique maximum,*

$$(4) \quad \kappa(s) < \kappa(s_0) =: k_{\max} , \text{ for all } s \neq s_0 ,$$

and that the maximum is non-degenerate, i.e.

$$(5) \quad k_2 := -\kappa''(s_0) \neq 0 .$$

Then, for all $n \in \mathbb{N}^$, there exists a sequence $\{\zeta_j^{(n)}\}_{j=1}^\infty \subset \mathbb{R}$ (which can be calculated recursively to any order) such that the n -th eigenvalue of \mathcal{H} $\mu^{(n)}(h)$ admits the following asymptotic expansion, when $h \searrow 0$,*

$$(6) \quad \mu^{(n)}(h) \sim \Theta_0 h - C_1 k_{\max} h^{3/2} + C_1 \Theta_0^{1/4} \sqrt{\frac{3k_2}{2}} (2n-1) h^{7/4} + h^{15/8} \sum_{j=0}^\infty h^{j/8} \zeta_j^{(n)} .$$

Remark 2.

Previous results on the bottom of the spectrum of $\mathcal{H}_{h,\Omega}$ were obtained in [HeMo], who gave the two first terms in the expansion of $\mu^{(1)}(h)$ (see [HeMo, Theorems 10.3 and 11.1]):

$$(7) \quad \mu^{(1)}(h) = \Theta_0 h - k_{\max} C_1 h^{3/2} + \mathcal{O}(h^{5/3}) .$$

Remark 3. *If the uniqueness condition in (4) is replaced by the assumption that there is a finite number of maxima (for which (5) is assumed to hold), we expect the existence of sequences of eigenvalues $z^{(n)}(h)$ corresponding to each maximum.*

For applications to bifurcations from the normal state in superconductivity it seems important to calculate the splitting between the ground and first excited states of $\mathcal{H}(h)$. Let us define

$$(8) \quad \Delta(h) = \mu^{(2)}(h) - \mu^{(1)}(h) .$$

Corollary 4.

Under the hypothesis of the theorem, $\Delta(h)$ admits the following asymptotics :

$$(9) \quad \Delta(h) \sim C_1 \Theta_0^{1/4} \sqrt{6k_2} h^{7/4} + h^{15/8} \sum_{j=0}^{\infty} h^{j/8} \xi_j .$$

where $\xi_j = \zeta_j^{(2)} - \zeta_j^{(1)}$.

The case where Ω is a disc has been analyzed in great detail in [BaPhTa], using the radial symmetry to reduce the problem to ordinary differential equations. In this case the splitting $\Delta(h)$ turns out to become zero for a sequence of values of h tending to 0. This is a complication in the analysis of bifurcation. Thus, in some sense, the more ‘generic’ situation considered here has a nicer property.

Thanks :

The authors are supported by the European Research Network ‘Postdoctoral Training Program in Mathematical Analysis of Large Quantum Systems’ with contract number HPRN-CT-2002-00277, and the ESF Scientific Programme in Spectral Theory and Partial Differential Equations (SPECT).

Address:

B. Helffer, Département de Mathématiques, UMR CNRS 8628,
Université Paris-Sud, Bât. 425, F-91405 Orsay Cedex (FRANCE).

REFERENCES

- [Ag] S. Agmon : *Lectures on exponential decay of solutions of second order elliptic equations.* Math. Notes, T. 29, Princeton University Press (1982).
- [BaPhTa] P. Bauman, D. Phillips, and Q. Tang : Stable nucleation for the Ginzburg-Landau system with an applied magnetic field. Arch. Rational Mech. Anal. 142, p. 1-43 (1998).
- [BeSt] A. Bernoff and P. Sternberg : Onset of superconductivity in decreasing fields for general domains. J. Math. Phys. 39, p. 1272-1284 (1998).
- [Bon] V. Bonnaillie : On the fundamental state for a Schrödinger operator with magnetic fields in domains with corners. Asymptotic Anal. 2004 (in Press).
- [BonDa] V. Bonnaillie and M. Dauge : Asymptotics for the fundamental state of the Schrödinger operator with magnetic field near a corner. In preparation, (2004).
- [FoHe1] S. Fournais and B. Helffer : Energy asymptotics for type II superconductors. Preprint 2004.
- [FoHe2] S. Fournais and B. Helffer : Accurate estimates for magnetic bottles in connection with superconductivity. Preprint 2004
- [Hel] B. Helffer : *Introduction to the semiclassical analysis for the Schrödinger operator and applications.* Springer lecture Notes in Math. 1336 (1988).
- [HeMo] B. Helffer and A. Morame : Magnetic bottles in connection with superconductivity. J. Funct. Anal. 185 (2), p. 604-680 (2001).
- [HeSj] B. Helffer and J. Sjöstrand : Multiple wells in the semiclassical limit I. Comm. Partial Differential Equations 9 (4), p. 337-408 (1984).
- [LuPa1] K. Lu and X-B. Pan : Estimates of the upper critical field for the Ginzburg-Landau equations of superconductivity. Physica D 127, p. 73-104 (1999).
- [LuPa2] K. Lu and X-B. Pan : Eigenvalue problems of Ginzburg-Landau operator in bounded domains. J. Math. Phys. 40 (6), p. 2647-2670, June 1999.

- [LuPa3] K. Lu and X-B. Pan : Gauge invariant eigenvalue problems on \mathbb{R}^2 and \mathbb{R}_+^2 . Trans. Amer. Math. Soc. 352 (3), p. 1247-1276 (2000).
- [PiFeSt] M. del Pino, P.L. Felmer, and P. Sternberg : Boundary concentration for eigenvalue problems related to the onset of superconductivity. Comm. Math. Phys. 210, p. 413-446 (2000).

Wave Front Set for Solutions to Schrödinger Equations

SHU NAKAMURA

In this talk, we discuss the wave front set for solutions to Schrödinger equation with variable coefficients. It is well-known that the propagation speed of the wave front set of solutions to Schrödinger equation is infinite, and hence we cannot expect the usual propagation theorem such as for the solutions to the wave equation. Instead, relations between the decay property of the initial condition and the wave front set of solutions have been studied, which is generally called (microlocal) smoothing properties. Here we present a different formulation, which is closer to the “propagation of singularity theorem”, at least in the spirit.

Part of results we discuss is joint work with André Martinez and Vania Sordani (Bologna University).

We consider a Schrödinger equation:

$$\frac{d}{dt}u(t) = -iHu(t), \quad u(0) = u_0 \in L^2(\mathbb{R}^d)$$

on $L^2(\mathbb{R}^d)$, where $d \geq 1$, and H is the Schrödinger operator defined by

$$H = \frac{1}{2} \sum_{i,j=1}^d D_j a_{jk}(x) D_k + V(x), \quad D_j = -i \frac{\partial}{\partial x_j}.$$

We suppose the coefficients $\{a_{ij}(x)\}$ and the potential $V(x)$ satisfy the following conditions:

Assumption A. $a_{ij}(x), V(x) \in C^\infty(\mathbb{R}^d; \mathbb{R})$ for $i, j = 1, \dots, d$, and there exist $\mu > 0$, and $C_\alpha > 0$ for each $\alpha \in \mathbb{Z}_+^d$ such that

$$|\partial_x^\alpha (a_{ij}(x) - \delta_{ij})| \leq C_\alpha \langle x \rangle^{-\mu - |\alpha|}, \quad |\partial_x^\alpha V(x)| \leq C_\alpha \langle x \rangle^{2 - \mu - |\alpha|},$$

for $x \in \mathbb{R}^d$. Moreover, H is elliptic, i.e., $\det(a_{ij}(x)) \neq 0$ for each $x \in \mathbb{R}^d$.

Then it is well-known that H is essentially self-adjoint on $C_0^\infty(\mathbb{R}^d)$. We denote the unique self-adjoint extension by the same symbol H . Thus, by the Stone theorem, $u(t) = e^{-itH}u_0$ is the solution to the Schrödinger equation with the initial condition $u(0) = u_0$. We consider the following quite basic PDE-type problem:

Problem: Describe the singularity of $u(t)$ in terms of u_0 .

We use the notion of the *wave front set* to describe singularity of solutions to the Schrödinger equation (see [6] Section X.10, or [8] Section VI.1 for the definition). We denote the wave front set of $u \in \mathcal{D}'(\mathbb{R}^d)$ by $WF(u) \subset \mathbb{R}^{2d}$.

Let us recall the propagation of singularity theorem for the wave equation. It shows that the propagation of the wave front set for the solutions to the wave equation is described by the geometric optics. We note the analogue of the geometric optics for Schrödinger equation is the classical mechanics, and the relationship is given by the WKB analysis. However, the WKB theory describe the semiclassical behavior of the solution, and it does not give any information about the singularity of solutions, at least directly. As we will see, the *high energy* classical mechanics gives us the description of the singularity of the solutions, and it is closely related to the scattering theory of the classical mechanics.

We denote the symbol of the kinetic energy part by $p(x, \xi)$, i.e.,

$$p(x, \xi) = \frac{1}{2} \sum_{i,j=1}^d a_{ij}(x) \xi_i \xi_j, \quad x, \xi \in \mathbb{R}^d.$$

We denote the solution to the Hamilton equation:

$$\frac{d}{dt}y(t) = \frac{\partial p}{\partial \xi}(y(t), \eta(t)), \quad \frac{d}{dt}\eta(t) = -\frac{\partial p}{\partial x}(y(t), \eta(t))$$

with initial condition $y(0) = x$, $\eta(0) = \xi$ by $y(t; x, \xi)$ and $\eta(t; x, \xi)$.

Definition 1: $(x, \xi) \in \mathbb{R}^{2d}$ is said to be *backward nontrapping* if $|y(t; x, \xi)| \rightarrow +\infty$ as $t \rightarrow -\infty$.

We say H is a *short-range perturbation* of $H_0 = -\frac{1}{2}\Delta$ (or simply short-range type) if Assumption A is satisfied with $\mu > 1$. In this case, if (x, ξ) is backward nontrapping, then it is well-known that there exists $(x_-, \xi_-) \in \mathbb{R}^{2d}$ such that

$$|y(t; x, \xi) - (x_- + t\xi_-)| \rightarrow 0 \quad \text{as } t \rightarrow -\infty.$$

Namely, the classical trajectory $y(t; x, \xi)$ approaches to a free motion $x_- + t\xi_-$ as $t \rightarrow -\infty$. The map:

$$S : (x, \xi) \mapsto (x_-, \xi_-)$$

is the classical (inverse) wave operator.

Theorem 1 ([5]) *Suppose Assumption A with $\mu > 1$, and suppose $(x, \xi) \in \mathbb{R}^{2d}$ is backward nontrapping. Let $u(t) = e^{-itH}u_0$ with $u_0 \in L^2(\mathbb{R}^d)$, and let $t > 0$. Then*

$$(x, \xi) \in WF(u(t)) \iff (x_-, \xi_-) \in WF(e^{-itH_0}u_0).$$

If the metric is flat, i.e., if $H = -\frac{1}{2}\Delta + V(x)$, then Theorem 1 implies that $WF(u(t)) = WF(e^{-itH_0}u_0)$. This observation suggests that $e^{itH_0}e^{-itH}$ is a pseudo-differential operator, and in fact we can prove it. This result and its generalization to non-flat case will be discussed in a forthcoming paper.

Recently, Hassel and Wunsch [2] have obtained different characterization of the wave front set of solutions to Schrödinger equations using the *quadratic scattering wave front set*. The setting and the formulation are quite different, and the relationship is not clear to the author.

If the perturbation is long-range type, i.e., if $0 < \mu \leq 1$, then the above theorem does not hold in general, and we need to replace the free propagator e^{-itH_0} by a different Fourier multiplier, quite similar to one appearing in the long-range scattering theory. This part is still in progress, and we do not discuss here. We have somewhat weaker result, which we discuss in the following. We recall that the classical motion not necessarily approaches to a free motion, but the asymptotic momentum $\xi_- := \lim_{t \rightarrow -\infty} \eta(t; x, \xi)$ does exist if the trajectory is nontrapping. We introduce the following notion of the wave front set:

Definition 2: Let $u \in \mathcal{S}'(\mathbb{R}^d)$. We say $(x, \xi) \in \mathbb{R}^{2d} \setminus 0$ is not in the *homogeneous wave front set* of u if there exists $a \in C_0^\infty(\mathbb{R}^{2d})$ such that $a(x, \xi) \neq 0$ and $\|a(hx, hD_x)u\|_{L^2} = O(h^N)$ as $h \rightarrow +0$ with any N . We denote $(x, \xi) \notin HWF(u)$ if this condition is satisfied, and denote the complement by $HWF(u)$.

Theorem 2 ([4]) *Suppose Assumption A with $\mu > 0$, and suppose $(x, \xi) \in \mathbb{R}^{2d}$ is backward nontrapping. Let $t_0 > 0$. If $(-t_0\xi_-, \xi_-) \notin HWF(u_0)$, then $(x_0, \xi_0) \notin WF(u(t_0))$.*

The microlocal smoothing property of Craig, Kappeler and Strauss [1] follows from Theorem 2. (In fact our result is more general, since they considered short-range case only.) It is also related to a work by Wunsch [9], though the relationship is not clear to the author. A similar theorem also holds for the analytic wave front set under the assumption of the analyticity of the coefficients. This result is proved by a recent joint work with Martinez and Sordani [3], and it is a generalization of results by Robbiano and Zuily [7].

REFERENCES

- [1] Craig, W., Kappeler, T., Strauss, W.: Microlocal dispersive smoothing for the Schrödinger equation, *Comm. Pure Appl. Math.* **48** (1996), 769–860.
- [2] Hassel, A., Wunsch, J.: The Schrödinger propagator for scattering metrics, Preprint 2003.
- [3] Martinez, A., Nakamura, S., Sordani, V.: Analytic smoothing effect for the Schrödinger equation with long-range perturbation. Preprint 2004, December. (available at http://www.ma.utexas.edu/mp_arc/)
- [4] Nakamura, S.: Propagation of the homogeneous wave front set for Schrödinger equations. To appear in *Duke Math. J.* **126** (2005)
- [5] Nakamura, S.: Wave front set for solutions to Schrödinger equation. Preprint 2004, Feb. (available at <http://www.ms.u-tokyo.ac.jp/~shu/>)
- [6] Read, M., Simon, B.: *The Methods of Modern Mathematical Physics. Vol. II.* Academic Press, 1975.
- [7] Robbiano, L., Zuily, C.: Microlocal analytic smoothing effect for Schrödinger equation, *Duke Math. J.* **100** (1999), 93–129.
- [8] Taylor, M. E.: *Pseudodifferential Operators*, Princeton Math. Series 34, 1981.
- [9] Wunsch, J.: Propagation of singularities and growth for Schrödinger operators, *Duke Math. J.* **98** (1999), 137–186.

Spectral shift function for self-adjoint operators without spectral gaps

D. R. YAFAEV

The concept of the spectral shift function first appeared in the work of I. M. Lifshits [7] in connection with the quantum theory of crystals. A mathematical theory of the SSF was soon thereafter constructed by M. G. Kreĭn in [5]. One of his results can be formulated in the following way. Let H_0 and H be self-adjoint operators with a trace class (denoted \mathfrak{S}_1) difference $V = H - H_0$. Then there exists a function $\xi(\lambda) = \xi(\lambda; H, H_0)$, $\xi \in L_1(\mathbb{R})$, known as the spectral shift function such that the trace formula

$$(1) \quad \text{Tr}\left(f(H) - f(H_0)\right) = \int_{-\infty}^{\infty} \xi(\lambda) f'(\lambda) d\lambda, \quad \xi(\lambda) = \xi(\lambda; H, H_0),$$

holds at least for all functions $f \in C_0^\infty(\mathbb{R})$. A relatively detailed presentation of the theory of the SSF can be found in [3] or [8].

If the operators H_0 and H have a common spectral gap, then the trace formula automatically remains true for a much wider class of the operators V . If, for instance, $\lambda = 0$ is a common regular point of the operators H_0 and H and $H^{-m} - H_0^{-m} \in \mathfrak{S}_1$ for some integer odd m , then the trace formula (1) for the pair H_0, H can be deduced from that for the pair H_0^{-m}, H^{-m} (see [8], for details).

A connection between scattering theory and the theory of the SSF was found by M. Sh. Birman and M. G. Kreĭn in [1]. Actually, they showed that the scattering matrix $S(\lambda; H, H_0)$ minus the identity operator I belongs to the trace class and

$$(2) \quad \det S(\lambda; H, H_0) = e^{-2\pi i \xi(\lambda; H, H_0)}$$

for almost all λ (from the core of the spectrum of the operator H_0).

Our goal is to extend the theory of the spectral shift function to the case where only the difference of some powers of the resolvents of self-adjoint operators belongs to the trace class. The main result is given by the following

Theorem 1. *Let, for a pair of self-adjoint operators H_0 and H , the assumption*

$$(H - z)^{-m} - (H_0 - z)^{-m} \in \mathfrak{S}_1$$

hold for some positive odd integer m and all z with $\Im z \neq 0$. Suppose that a function $f(\lambda)$ has two bounded derivatives and

$$\partial^\alpha (f(\lambda) - f_0 \lambda^{-m}) = O(|\lambda|^{-m-\epsilon-\alpha}), \quad \alpha = 0, 1, 2, \quad \epsilon > 0,$$

where the constant f_0 is the same for $\lambda \rightarrow \infty$ and $\lambda \rightarrow -\infty$. Then

$$f(H) - f(H_0) \in \mathfrak{S}_1$$

and there exists a function $\xi(\lambda; H, H_0)$ satisfying the condition

$$\int_{-\infty}^{\infty} |\xi(\lambda; H, H_0)| (1 + |\lambda|)^{-m-1} d\lambda < \infty$$

such that the trace formula (1) is true. Moreover, for the corresponding scattering matrix $S(\lambda; H, H_0)$, the operator $S(\lambda; H, H_0) - I \in \mathfrak{S}_1$ and the relation (2) holds for almost all λ .

Note that in the case $m = 1$ Theorem 1 reduces to a well-known result of M. G. Kreĭn [6]. Somewhat different general conditions for the existence of the spectral shift function were given by L. S. Koplienko [4].

Our proof of Theorem 1 relies on its reduction to the special case $m = 1$ with the help of the theory of Double Operator Integrals developed by M. Sh. Birman and M. Z. Solomyak (see, e.g., [2]).

REFERENCES

- [1] M. Sh. Birman and M. G. Kreĭn, *On the theory of wave operators and scattering operators*, Soviet Math. Dokl. **3** (1962), 740–744.
- [2] M. Sh. Birman and M. Z. Solomyak, *Double operator integrals in a Hilbert space*, Int. Eq. Oper. Th. **47** (2003), 131–168.
- [3] M. Sh. Birman and D. R. Yafaev, *The spectral shift function. The papers of M. G. Kreĭn and their further development*, St. Petesburg Math. J. **4** (1993), 833–870.
- [4] L. S. Koplienko, *Local criteria for the existence of the spectral shift function*, Zap. Nauchn. Sem. LOMI Steklov **73** (1977), 102–117. (Russian)
- [5] M. G. Kreĭn, *On the trace formula in perturbation theory*, Mat. Sb. **33** (1953), 597–626. (Russian)
- [6] M. G. Kreĭn, *On perturbation determinants and the trace formula for unitary and selfadjoint operators*, Soviet Math. Dokl. **3** (1962) 707–710.
- [7] I. M. Lifshits, *On a problem in perturbation theory*, Uspehi Mat. Nauk **7** (1952) 171–180. (Russian)
- [8] D. R. Yafaev, *Mathematical scattering theory*, Amer. Math. Soc., Providence, Rhode Island, 1992.

Uniform Magnetic Lieb-Thirring inequalities

LÁSZLÓ ERDŐS

(joint work with Jan Philip Solovej)

Lieb-Thirring inequalities refer to estimates that bound moments of negative eigenvalues of Schrödinger type operators in terms of the external fields. They play a fundamental role in various results concerning localized many-fermion systems. Most notably, the ground state energy of the many-body Hamiltonian in many cases is related to the sum of the negative eigenvalues of an effective one-body Hamiltonian. Among other useful applications, Lieb-Thirring inequalities stand behind the most effective and elegant proofs of stability of matter. They also serve as a powerful a priori estimate for the semiclassical analysis of the many-fermion ground state.

We focus on the particular case of *magnetic Lieb-Thirring (MLT) inequalities*. They estimate moments of negative eigenvalues $e_1(H) \leq e_2(H) \leq \dots \leq 0$ of the Pauli operator

$$(1) \quad H := [\boldsymbol{\sigma} \cdot (-i\nabla + \mathbf{A})]^2 + V$$

on $L^2(\mathbf{R}^3, \mathbf{C}^2)$ with a vector potential \mathbf{A} , magnetic field $\mathbf{B} := \nabla \times \mathbf{A}$ and external potential V . Here $\boldsymbol{\sigma} \cdot \mathbf{v} = \sigma^1 v_1 + \sigma^2 v_2 + \sigma^3 v_3$, $\mathbf{v} \in \mathbf{R}^3$, and $\sigma^1, \sigma^2, \sigma^3$ are the Pauli matrices. Unlike in the nonmagnetic case, where the *optimal form* of the

estimates is well-established and the remaining main challenge is to find the *optimal constants*, the magnetic case is much less understood.

For a constant magnetic field \mathbf{B} , the inequality

$$(2) \quad \sum_j |e_j(H)| \leq (\text{const}) \left(\int_{\mathbf{R}^3} |\mathbf{B}| [V]_-^{3/2} + \int_{\mathbf{R}^3} [V]_-^{5/2} \right),$$

proven in [LSY], is optimal, apart from the constant, where $[a]_- := -\min\{0, a\}$ denotes the negative part of a . It has seemed to be reasonable to conjecture that (2) also holds for an arbitrary magnetic field. However, such a naive generalization fails for two reasons.

Firstly, even when \mathbf{B} has constant direction in \mathbf{R}^3 , (2) can be correct only if $|\mathbf{B}(x)|$ is replaced by an effective field strength, $B^*(x)$, obtained by averaging $|\mathbf{B}|$ locally on the magnetic lengthscale, $|\mathbf{B}|^{-1/2}$.

Secondly, the existence of the celebrated Loss-Yau zero modes [LY] contradicts (2). Indeed, for certain magnetic fields with nonconstant direction the Dirac operator $\mathcal{D} := \boldsymbol{\sigma} \cdot (-i\nabla + \mathbf{A})$ has a nontrivial L^2 -kernel. In this case a small potential perturbation of \mathcal{D}^2 shows that $\sum_j |e_j(H)|$ behaves as $\int n(x)[V(x)]_- dx$, i.e. it is linear in $[V]_-$. Here $n(x)$ is the density of zero modes, $n(x) = \sum_j |u_j(x)|^2$, where $\{u_j\}$ is an orthonormal basis in $\text{Ker } \mathcal{D}$. Thus an extra term linear in $[V]_-$ must be added to (2) and $n(x)$ has to be estimated.

Let

$$H(h, b) := [\boldsymbol{\sigma} \cdot (-ih\nabla + b\mathbf{A})]^2 + V$$

be the Pauli operator with the semiclassical parameter h and a field strength parameter b . The semiclassical formula for the sum of the negative eigenvalues, i.e. the asymptotic formula for $\sum_j |e_j(H(h, b))|$ as $h \rightarrow 0$, behaves linearly in the field strength for a constant magnetic field ([LSY]). This fact suggests that $\sum_j |e_j(H)|$ may be bounded by an expression that grows only with the first power of $|\mathbf{B}|$ even for nonconstant magnetic fields.

Our goal is to establish such MLT estimates with as few technical assumptions on \mathbf{B} as possible and no technical assumptions on V . The density $n(x)$ has a dimension $(\text{length})^{-3}$. Since $|\mathbf{B}|$ has dimension $(\text{length})^{-2}$, we need to introduce an extra lengthscale to be able to bound $n(x)$ by the magnetic field. We will therefore make assumptions on certain derivatives of the field.

In almost all previous Lieb-Thirring bounds, the density $n(x)$ was estimated by a function that behaves quantitatively as $|\mathbf{B}(x)|^{3/2}$. This power was reduced to 5/4 in [ES-I] with a further unnatural $V \in W^{1,1}$ assumption on the potential. A worse power, 17/12, was obtained in [BFG] but without further assumptions on the potential.

For our theorem we assume that $\mathbf{B}(x) \neq 0$ for all $x \in \mathbf{R}^3$, i.e. the unit vectorfield $\mathbf{n} := \mathbf{B}/|\mathbf{B}|$ is well defined. We also assume that the vectorfield \mathbf{n} satisfies the following global regularity condition

$$(3) \quad L_{\mathbf{n}}^{-1} := \sum_{\gamma=1}^5 \|\nabla^\gamma \mathbf{n}\|_\infty^{1/\gamma} < \infty.$$

For any $L > 0$, $x \in \mathbf{R}^3$ we also define

$$B_L^*(x) := \sup\{|\mathbf{B}(y)| : |y - x| \leq L\} + L \cdot \sup\{|\nabla \mathbf{B}(y)| : |y - x| \leq L\}.$$

Theorem 1 (Magnetic Lieb-Thirring inequality). [ES-IV] *For any $0 < L \leq L_{\mathbf{n}}$, the sum of the negative eigenvalues, $e_1(H) \leq e_2(H) \leq \dots \leq 0$, of H satisfies*

$$(4) \quad \sum_j |e_j(H)| \leq (\text{const}) \left(L^{-1} \int (B_L^* + L^{-2}) [V]_- + \int B_L^* [V]_-^{3/2} + \int [V]_-^{5/2} \right).$$

The density of Loss-Yau zero modes is estimated by

$$n(x) \leq L^{-1} (B_L^*(x) + L^{-2}).$$

The linear power of $|\mathbf{B}|$ in the estimate reflects the basic fact that the space with a magnetic field cannot be considered isotropic: the quantum motion parallel with the magnetic field behaves differently than the transversal one. The magnetic field affects only the two-dimensional transversal motion. All MLT estimates that yield $|\mathbf{B}|^{3/2}$ behaviour neglect this geometric fact by simply comparing the magnetic problem with a three dimensional nonmagnetic one, usually via a diamagnetic inequality that loses the anisotropic feature of the problem. The typical estimate is of the form

$$(5) \quad \mathcal{D}^2 \geq b^{-1} \mathcal{D}^2 = b^{-1} [(-i\nabla + \mathbf{A})^2 + \boldsymbol{\sigma} \cdot \mathbf{B}],$$

where $b := \|\mathbf{B}\| \gg 1$ is some (local) norm of \mathbf{B} . The kinetic energy is scaled down so that the dangerous $\boldsymbol{\sigma} \cdot \mathbf{B}$ becomes bounded uniformly in b . The magnetic Laplacian can then be controlled by the nonmagnetic Laplacian, $-\Delta$, but the factor b^{-1} now affects all three coordinates, yielding a scaling of $b^{3/2}$. The key to our theorem is to separate the parallel and transversal motions and use a crude estimate similar to (5) only in the two-dimensional transversal kinetic energy.

Our theorem uses only natural assumptions on V and it gives the correct (linear) dependence on the field strength $|\mathbf{B}|$. However, the original magnetic field \mathbf{B} is replaced by an effective field $B_L^* + L^{-2}$ that involves the global L^∞ -norm of \mathbf{n} . In particular the estimate (4) is sensitive to the behavior of the magnetic field far away from the support of $[V]_-$. Hence the irregular behaviour of \mathbf{n} far away from the support of $[V]_-$ renders our estimate large despite that it should not substantially influence the negative spectrum.

In a separate work [ES-III] we also proved a magnetic Lieb-Thirring bound that enjoys a *locality property*. More precisely, the constant $L_{\mathbf{n}}$ was replaced by a function $L(x)$ describing the local variation lengthscale of the magnetic field. The precise definition is somewhat complicated, but it depends only locally on \mathbf{B} . In particular the inverse lengthscale $L^{-1}(x)$ can be bounded by $c\delta^{-1}$ if \mathbf{B} vanishes in a δ -neighborhood of x .

The proof of this second theorem is much more involved. The complications are due to the lack of effective offdiagonal bounds on the resolvent of the Pauli operator, $(\mathcal{D}^2 + E)^{-1}(x, y)$. For constant magnetic field, the resolvent decays on the magnetic lengthscale $B^{-1/2}$ in the direction perpendicular to the field:

$$(\mathcal{D}^2 + E)^{-1}(x, y) \sim e^{-cB(x_\perp - y_\perp)^2}$$

but similar estimate is unknown for a general field. This problem is closely related to the poorly understood structure of the Loss-Yau zero modes.

It is amusing to note that it was a substantial endeavour to show that a zero mode may exist at all [LY], and that multiple zero modes may also occur [ES-II]. On the other hand, it seems also quite difficult to give an upper bound on their number in terms of the first power of the field strength.

REFERENCES

- [BFG] L. Bugliaro, C. Fefferman and G. M. Graf: *A Lieb-Thirring bound for a magnetic Pauli Hamiltonian, II*, Rev. Mat. Iberoamericana, **15**, 593-619 (1999)
- [ES-I] L. Erdős and J. P. Solovej: *Semiclassical eigenvalue estimates for the Pauli operator with strong non-homogeneous magnetic fields. I. Non-asymptotic Lieb-Thirring type estimate*. Duke J. Math. **96**, 127-173 (1999)
- [ES-II] L. Erdős and J. P. Solovej: *The kernel of Dirac operators on S^3 and \mathbf{R}^3* . Rev. Math. Phys. **13** No. 10, 1247-1280 (2001)
- [ES-III] L. Erdős and J.P. Solovej, *Uniform Lieb-Thirring inequality for the three dimensional Pauli operator with a strong non-homogeneous magnetic field*. Ann. Inst. H. Poincaré **5**, 671-741 (2004)
- [ES-IV] L. Erdős and J.P. Solovej, *Magnetic Lieb-Thirring inequalities with optimal dependence on the field strength*. J. Stat. Phys. **116**, 475-506 (2004)
- [LSY] E. H. Lieb, J. P. Solovej and J. Yngvason: *Asymptotics of heavy atoms in high magnetic fields: II. Semiclassical regions*. Commun. Math. Phys. **161**, 77-124 (1994)
- [LY] M. Loss and H.-T. Yau: *Stability of Coulomb systems with magnetic fields: III. Zero energy bound states of the Pauli operator*. Commun. Math. Phys. **104**, 283-290 (1986)

Old and New Tales about Lifshitz Tails

WERNER KIRSCH

We consider random Schrödinger operators $H_\omega = H_0 + V_\omega$ with V_ω either of alloy type or of Poisson type.

By alloy type we mean potentials of the form

$$(1) \quad V_\omega(x) = \sum_{i \in \mathbb{Z}^d} q_i(\omega) f(x - i)$$

with independent identically distributed random variables q_i . For the Poisson model the potential V_ω is given by

$$(2) \quad V_\omega(x) = \sum f(x - \xi_i(\omega))$$

where the $\{\xi_i\}$ are Poisson distributed random points.

In both cases the function f , also called the single site potential, has to decay fast enough at infinity to ensure convergence of the series (1), e.g.

$$(3) \quad |f(x)| \leq \frac{c}{1 + |x|^\alpha}$$

with $\alpha > d$, $|x|$ large.

We study the integrated density of states $N(E)$ for these operators. $N(E)$ is defined as a thermodynamic limit in the following way: Let Λ_L be the cube of side length L around the origin and restrict H_ω to $\ell^2(\Lambda_L)$ with appropriate boundary conditions (Dirichlet, say). The corresponding operator H_L has a purely discrete spectrum. Let us denote its eigenvalues by $E_1(H_L) \leq E_2(H_L) \leq \dots$, repeated according to multiplicity.

For any E we set $N_L(E) = \frac{1}{L^d} \#\{E_n(H_L) \leq E\}$.

Under very weak conditions on V_ω it is known that N_L converges (for all but countably many E at least) to a nonrandom limit $N(E)$, the integrated density of states.

The physicist Lifshitz [1] observed in 1964 that the low energy behavior of $N(E)$ of random potentials is drastically different from the one for periodic potentials. In fact, Lifshitz argued that in the ordered (i.e. periodic) case

$$(4) \quad N(E) \sim (E - E_0)^{\frac{d}{2}}$$

as $E \searrow E_0 = \inf \sigma(H_{per})$. For random operators Lifshitz obtained

$$(5) \quad N(E) \sim e^{-c(E-E_0)^{-\frac{d}{2}}}$$

as $E \searrow E_0$. This super exponential decay is nowadays called Lifshitz tail behavior.

Lifshitz' arguments for his results were convincing but not mathematically rigorous.

The first mathematical proof of (5) was given by Donsker and Varadhan in [2].

Their proof relies on the machinery of the Donsker-Varadhan large deviations results. For their proof to work Donsker and Varadhan need that the single-site potential decays fast enough, namely:

$$(6) \quad |f(x)| \leq \frac{c}{1 + |x|^\alpha}$$

with $\alpha > d + 2$.

Pastur [3] proved that for slower decay, i.e.

$$(7) \quad f(x) \sim \frac{c}{1 + |x|^\alpha}$$

with $d < \alpha < d + 2$ the Lifshitz behavior (5) is changed to

$$(8) \quad N(E) \sim e^{-c(E-E_0)^{-\frac{d}{\alpha-d}}}.$$

In the eighties the so called Dirichlet-Neumann-bracketing technique was used to prove Lifshitz tails (as in (5) or in (8)) for a greater variety of random potentials

([4], [5], [6]). This technique is much simpler than the Donsker-Varadhan method. It is also much closer to the original physical arguments by Lifshitz.

Recently, in [7] single-site potentials with anisotropic decay were considered.

Suppose that $x = (x_1, x_2)$ $x_1 \in \mathbb{R}^{d_1}$, $x_2 \in \mathbb{R}^{d_2}$ and

$$(9) \quad |f(x)| \sim \frac{c}{1 + |x_1|^{\alpha_1} + |x_2|^{\alpha_2}}.$$

We define $\gamma_i = \frac{d_i}{d_k}$ and $\gamma = \gamma_1 + \gamma_2$. Then

$$(10) \quad N(E) \sim e^{-c(E-E_0)^{-\eta}}$$

where

$$(11) \quad \eta = \max\left(\frac{d_1}{2}, \frac{\gamma_1}{1-\gamma}\right) + \max\left(\frac{d_2}{2}, \frac{\gamma_2}{1-\gamma}\right).$$

If we introduce a constant magnetic field into the Hamiltonian the Lifshitz behavior is qualitatively changed.

For example for $d = 2$ and $f(x) \sim \frac{c}{1+|x|^\alpha}$ it was proved [8] that

$$(12) \quad N(E) \sim e^{-c(E-E_0)^{-\frac{d}{\alpha-d}}}.$$

even if $\alpha > d + 2$. L. Erdős [9] proved that for compactly supported f $N(E)$ decays polynomially. For $d = 3$ see ([10], [11]).

Finally, we mention that Lifshitz tails may also exist at internal band edges ([12], [13]) as well as for random surface potentials ([14], [15]).

REFERENCES

- [1] I. M. Lifshitz. Structure of the energy spectrum of the impurity bands in disordered solid solutions. *Sov. Phys. JETP*, 17:1159–1170, 1963. Russian original: *Zh. Eksp. Ter. Fiz.*, 44:1723–1741, 1963.
- [2] M.D. Donsker, S.R.S. Varadhan. Asymptotic for a Wiener sausage. *Commun. Pure Appl. Math.* 28, 525–565 (1975)
- [3] L. A Pastur. Behavior of some Wiener integrals as $t \rightarrow \infty$ and the density of states of Schrödinger equations with random potential. *Theor. Math. Phys.*, 32:615–620, 1977. Russian original: *Teor. Mat. Fiz.*, 6:88–95, 1977.
- [4] W. Kirsch and F. Martinelli. Large deviations and Lifshitz singularity of the integrated density of states of random Hamiltonians. *Commun. Math. Phys.*, 89:27–40, 1983.
- [5] W. Kirsch and B. Simon. Lifshitz tails for periodic plus random potentials. *J. Stat. Phys.*, 42:799–808, 1986.
- [6] G. A. Mezincescu. Lifshitz singularities for periodic operators plus random potential. *J. Stat. Phys.*, 49:1181–1190, 1987.
- [7] W. Kirsch; S. Warzel. Lifshitz tails caused by anisotropic decay: the emergence of a quantum-classical regime. Preprint math-ph/0310033
- [8] K. Broderix, D. Hundertmark, W. Kirsch, and H. Leschke. The fate of Lifshitz tails in magnetic fields. *J. Stat. Phys.*, 80:1–22, 1995.
- [9] L. Erdős. Lifshitz tail in a magnetic field: the nonclassical regime. *Probab. Theory Relat. Fields*, 112:321–371, 1998.

- [10] S. Warzel. *On Lifshits tails in magnetic fields*. Logos, Berlin, 2001. PhD thesis, University Erlangen-Nürnberg 2001.
- [11] H. Leschke, S. Warzel. Quantum-classical transitions in Lifshitz tails with magnetic fields. *Physical Review Letters* 92, 086402: 1-4 (2004)
- [12] F. Klopp. Internal Lifshits tails for random perturbations of periodic Schrödinger operators. *Duke Math. J.*, 98:335–369, 1999. Erratum: mp-arc 00-389.
- [13] F. Kopp, T. Wolff. Lifshitz tails for 2-dimensional random Schrödinger operators. *J. Anal Math.*, 88: 63-147 (2002). Dedicated to the memory of Tom Wolff
- [14] W. Kirsch, F. Klopp. The band edge behavior of the density of surfacic states. to appear in MPAG, math-ph/0407051
- [15] W. Kirsch, S. Warzel. Anderson localization and Lifshitz tails for random surface potentials. Preprint math-ph/0412079

Liouville theorems on abelian coverings of compact manifolds

PETER KUCHMENT

(joint work with Yehuda Pinchover, Technion, Israel)

The classical Liouville theorem claims that any harmonic function in \mathbb{R}^n of a polynomial growth is in fact a polynomial. In particular, the space of all harmonic functions that grow not faster than $C(1+|x|)^N$, is of finite dimension $\binom{n+N}{N} - \binom{n+N-2}{N-2}$. Analogously, the space of holomorphic function in \mathbb{C}^n of same growth consists of holomorphic polynomials of order N . The problem of extending this result to more general elliptic operators and/or to Laplace-Beltrami operators on general Riemannian manifolds of non-negative Ricci curvature was suggested in work of S. T. Yau [14]. One is interested in finite dimensionality of the spaces of solutions of a prescribed polynomial growth, estimates of (or even formulas for) their dimensions, and structure of these solutions (see [3, 7, 8] and references therein). Yau's conjecture on validity of the Liouville theorem for Riemannian manifolds of non-negative Ricci curvature was proven in full generality in [3] (see a description of previous partial results in [7, 8]).

An amazing case was discovered in [1, 13], where divergence form periodic elliptic equations $Lu = - \sum_{1 \leq i, j \leq n} (a^{i,j}(x)u_{x_i})_{x_j} = 0$ were considered. It was shown that the space of solutions of polynomial growth of order at most N of $Lu = 0$ has the same dimension as the space of harmonic polynomials of the same rate of growth. Any such solution is representable in the Floquet form $v(x) = \sum_{j=(j_1, \dots, j_n) \in \mathbb{Z}_+^n} x^j p_j(x)$, with periodic coefficients $p_j(x)$.

The natural questions to ask are: Is it important that the operator is of divergence form? What can be said about more general periodic (including higher order and matrix) equations? Is it possible to determine for a given periodic elliptic equation whether the Liouville theorem holds? How crucial is the usage in [1, 13] of homogenization theory tools (which automatically restrict the class

of equations)? Can these results be generalized for abelian coverings of compact manifolds?

Some partial answers to these questions were obtained in simultaneous papers [5, 9]. In [9], the results were generalized to second order periodic operators without lower order terms. At the same time, [5] contained a necessary and sufficient condition for the validity of the Liouville theorem for a periodic elliptic operator in \mathbb{R}^n , as well as (in most cases implicit) description of the dimensions of the corresponding spaces of solutions.

Simultaneously, an activity has existed of studying Liouville theorems for holomorphic functions on complex analytic manifolds (see e.g., [10, 11, 12]). In particular, one asks whether Liouville theorems for holomorphic functions hold for coverings of compact analytic manifolds. One of the results in [10] shows that the space of such bounded functions is finite-dimensional on nilpotent coverings of compact complex analytic manifolds. It was not clear whether one could say the same about the spaces of functions of given polynomial growth, except in the Kähler case [2].

The talk describes the results of [6] that clarify these issues for elliptic equations and systems (including overdetermined ones) on abelian coverings of compact Riemannian manifolds and holomorphic functions on abelian coverings of compact complex manifolds. The crucial techniques used come from the Floquet theory and are related to spectral notions common to the solid state physics.

Let $X \xrightarrow{G} M$ be an abelian covering of a compact d -dimensional Riemannian manifold M with an abelian deck group G (w.l.o.g. one can assume $G = \mathbb{Z}^n$). Let P be an elliptic G -periodic operator on X , with sufficiently smooth coefficients. For any character χ of G , one can consider a “twisted” version $P(\chi)$ of P on M that acts in sections of the linear bundle on M determined by χ (it is the push-down of P considered on χ -automorphic functions on X). In “normal” cases, the spectra of all $P(\chi)$ are discrete. The spectrum of $P(\chi)$ as a multiple-valued function of the character χ is called in solid state physics the *dispersion curve* or *dispersion relation* of P . The *Fermi surface* F of P is the set of unitary characters χ s.t. $Pu = 0$ has a non-zero χ -automorphic solution (i.e., F is the zero level set of the dispersion relation).

We say that the Liouville theorem holds to an order N for $Pu = 0$, if the space $V_N(P)$ of solutions of the equation with a bound $|u(x)| \leq C(1 + \rho(x))^N$ is finite dimensional. Here $\rho(x)$ is the distance of $x \in X$ from a fixed point $x_0 \in X$.

The theorem below describes our main results for the elliptic case.

Theorem 1

- (1) *If Liouville theorem for the equation $Pu = 0$ holds to an order $N \geq 0$, it holds to any order.*
- (2) *In order for the Liouville theorem to hold, it is necessary and sufficient that the Fermi surface F consists of finitely many points (this essentially means that one should expect the Liouville theorem to hold only when zero is at an edge of the spectrum of P).*

- (3) *If the Liouville theorem holds, then under some genericity condition on the operator P , the dimension of the space $V_N(P)$ can be computed explicitly in terms of the first non-zero term of the Taylor expansion of the dispersion curve near its zeros.*
- (4) *Under the same conditions, one can describe a constant coefficient ('homogenized') linear differential operator $\Lambda(D)$ on \mathbb{R}^n , such that there is a one-to-one correspondence between polynomial solutions of $\Lambda v = 0$ on \mathbb{R}^n and polynomially growing solutions of $Pu = 0$ on X .*

Similar results hold for overdetermined elliptic systems, including Cauchy-Riemann $\bar{\partial}$ operators. Here one obtains in particular the following

Theorem 2 *Let $X \rightarrow M$ be an abelian covering of a compact complex analytic manifold M and X be equipped with a periodic with respect to the deck group Riemannian metric. Then for any N the space of holomorphic functions on X of the polynomial growth of order N is finite dimensional. All such functions are polynomials of a fixed finite set of holomorphic functions.*

The proofs of the results dependent upon the techniques of Floquet theory [4]. This work was partially supported by NSF and BSF grants.

REFERENCES

- [1] M. Avellaneda and F.-H. Lin, *Un theoreme de Liouville pour des equations elliptiques a coefficients periodiques*, C. R. Acad. Sci. Paris, **309** (1989), 245–250.
- [2] A. Brudnyi, *Holomorphic functions of polynomial growth on abelian coverings of a compact complex manifold*, Comm. Anal. Geom. **6** (1998), no. 3, 485–510.
- [3] T. H. Colding and W. P. Minicozzi, *Harmonic functions on manifolds*, Ann. of Math. **146** (1997), 725–747.
- [4] P. Kuchment, *Floquet Theory for Partial Differential Equations*, Birkhäuser Verlag, Basel, 1993.
- [5] P. Kuchment and Y. Pinchover, *Integral representations and Liouville theorems for solutions of periodic elliptic equations*, J. Funct. Anal. **181**(2001), 402–446.
- [6] P. Kuchment and Y. Pinchover, *Liouville theorems and spectral edge behavior on abelian coverings of compact manifolds*, in preparation.
- [7] P. Li, *Curvature and function theory on Riemannian manifolds*, Surveys in Diff. Geom.: Papers dedicated to Atiyah, Bott, Hirzebruch, and Singer v. VII, International Press (2000), 375–432.
- [8] P. Li and J. P. Wang, *Counting dimensions of L-harmonic functions*, Ann. Math. **152** (2000), 645–658.
- [9] P. Li and J. Wang, *Polynomial growth solutions of uniformly elliptic operators of non-divergence form*, Proc. Amer. Math. Soc., **129** (2001), no. 12, 3691–3699.
- [10] V. Ya. Lin, *Liouville coverings of complex spaces, and amenable groups*, Math. USSR Sbornik, **60** (1998), no. 1, 197–216.
- [11] V. Ya. Lin and M. G. Zaidenberg, *Finiteness theorems for holomorphic maps*, in Encyclopaedia of Math. Sci., Vol.9. Several Complex Variables III. Berlin-Heidelberg-New York, Springer Verlag, 1989, 113–172.
- [12] V. Lin and M. Zaidenberg, *Liouville and Carathéodory coverings in Riemannian and complex geometry*, in Voronezh Winter Mathematical Schools, 111–130, Amer. Math. Soc. Transl. Ser. 2, 184, Amer. Math. Soc., Providence, RI, 1998.
- [13] J. Moser and M. Struwe, *On a Liouville-type theorem for linear and nonlinear elliptic differential equations on a torus*, Bol. Soc. Brasil. Mat. **23** (1992), 1–20.

- [14] S. T. Yau, *Harmonic functions on complete Riemannian manifolds*, Commun. Pure Appl. Math. **28** (1975), 201–228.

Homogenization problem for the stationary periodic Maxwell system

T. A. SUSLINA

1. We study the homogenization problem for the stationary periodic Maxwell system in the small period limit. There is a vast literature on this problem. In particular, it was discussed in the books [1,2]. However, known results provide only the weak convergence of solutions. We report on the new results [5,6] about approximations of the solutions in the $L_2(\mathbb{R}^3)$ -norm. The results are based on the abstract operator theory approach developed in [3,4].

2. Statement of the problem. Let Γ be a lattice in \mathbb{R}^3 , and let Ω be the cell of Γ . Suppose that the permittivity $\eta(\mathbf{x})$ and the permeability $\mu(\mathbf{x})$ are Γ -periodic measurable (3×3) -matrix-valued functions in \mathbb{R}^3 with real entries, and

$$c_0 \mathbf{1} \leq \eta(\mathbf{x}) \leq c_1 \mathbf{1}, \quad c_0 \mathbf{1} \leq \mu(\mathbf{x}) \leq c_1 \mathbf{1}, \quad \mathbf{x} \in \mathbb{R}^3, \quad 0 < c_0 \leq c_1 < \infty. \quad (1)$$

Here $\mathbf{1}$ is the identity matrix. We put $\mathfrak{G} = L_2(\mathbb{R}^3; \mathbb{C}^3)$. By $\mathfrak{G}(\eta^{-1}) = L_2(\mathbb{R}^3; \mathbb{C}^3; \eta^{-1})$ we denote the "weighted" space with the norm $(\eta^{-1} \mathbf{f}, \mathbf{f})_{\mathfrak{G}}^{1/2}$. The space $\mathfrak{G}(\mu^{-1}) = L_2(\mathbb{R}^3; \mathbb{C}^3; \mu^{-1})$ is defined in a similar way. We denote $J = \{\mathbf{f} \in \mathfrak{G} : \operatorname{div} \mathbf{f} = 0\}$. In what follows, \mathbf{u} and \mathbf{v} stand for the electric and magnetic field intensity, respectively, $\mathbf{w} = \eta \mathbf{u}$ is the electric displacement vector, and $\mathbf{z} = \mu \mathbf{v}$ is the magnetic induction vector. We represent the Maxwell operator $\mathcal{M} = \mathcal{M}(\eta, \mu)$ in terms of \mathbf{w} and \mathbf{z} assuming that they are divergence free. Then \mathcal{M} acts in the space $J \oplus J$ and is defined by the relations

$$\mathcal{M}(\eta, \mu) = \begin{pmatrix} 0 & i \operatorname{rot} \mu^{-1} \\ -i \operatorname{rot} \eta^{-1} & 0 \end{pmatrix}, \quad (2)$$

$$\operatorname{Dom} \mathcal{M}(\eta, \mu) = \{(\mathbf{w}, \mathbf{z}) : \mathbf{w} \in J, \mathbf{z} \in J, \operatorname{rot} \eta^{-1} \mathbf{w} \in \mathfrak{G}, \operatorname{rot} \mu^{-1} \mathbf{z} \in \mathfrak{G}\}.$$

The operator \mathcal{M} is selfadjoint in $J \oplus J$ treated as a subspace of $\mathfrak{G}(\eta^{-1}) \oplus \mathfrak{G}(\mu^{-1})$.

Let ε be a parameter. We denote $\eta^\varepsilon(\mathbf{x}) = \eta(\varepsilon^{-1} \mathbf{x})$, $\mu^\varepsilon(\mathbf{x}) = \mu(\varepsilon^{-1} \mathbf{x})$. Consider the family of operators $\mathcal{M}_\varepsilon = \mathcal{M}(\eta^\varepsilon, \mu^\varepsilon)$. Our goal is to study the behavior of the resolvent $(\mathcal{M}_\varepsilon - iI)^{-1}$ as $\varepsilon \rightarrow 0$. Consider the equation

$$(\mathcal{M}_\varepsilon - iI) \begin{pmatrix} \mathbf{w}_\varepsilon \\ \mathbf{z}_\varepsilon \end{pmatrix} = \begin{pmatrix} \mathbf{q} \\ \mathbf{r} \end{pmatrix}, \quad \mathbf{q}, \mathbf{r} \in J. \quad (3)$$

The corresponding intensities are given by $\mathbf{u}_\varepsilon = (\eta^\varepsilon)^{-1} \mathbf{w}_\varepsilon$ and $\mathbf{v}_\varepsilon = (\mu^\varepsilon)^{-1} \mathbf{z}_\varepsilon$. We are interested in the behavior of all four fields \mathbf{u}_ε , \mathbf{v}_ε , \mathbf{w}_ε , \mathbf{z}_ε as $\varepsilon \rightarrow 0$.

3. Results. It is useful to represent the solution components as the sums $\mathbf{w}_\varepsilon = \mathbf{w}_\varepsilon^{(q)} + \mathbf{w}_\varepsilon^{(r)}$, $\mathbf{z}_\varepsilon = \mathbf{z}_\varepsilon^{(q)} + \mathbf{z}_\varepsilon^{(r)}$, where the pair $\mathbf{w}_\varepsilon^{(q)}$, $\mathbf{z}_\varepsilon^{(q)}$ is the solution of (3) with $\mathbf{r} = 0$ and the pair $\mathbf{w}_\varepsilon^{(r)}$, $\mathbf{z}_\varepsilon^{(r)}$ is the solution of (3) with $\mathbf{q} = 0$. The fields \mathbf{u}_ε and \mathbf{v}_ε are represented in a similar way. For "half of the fields", namely, for $\mathbf{v}_\varepsilon^{(r)}$, $\mathbf{z}_\varepsilon^{(r)}$ and $\mathbf{u}_\varepsilon^{(q)}$, $\mathbf{w}_\varepsilon^{(q)}$ we obtain uniform approximations in the \mathfrak{G} -norm.

These approximations are of precise order with respect to parameter ε . For the remaining fields $\mathbf{v}_\varepsilon^{(q)}$, $\mathbf{z}_\varepsilon^{(q)}$ and $\mathbf{u}_\varepsilon^{(r)}$, $\mathbf{w}_\varepsilon^{(r)}$ we still have only weak convergence (which was known before).

Consider the case where $\mathbf{q} = 0$ in detail. Then equation (3) takes the form

$$\left. \begin{aligned} \mathbf{w}_\varepsilon^{(r)} &= \text{rot} (\mu^\varepsilon)^{-1} \mathbf{z}_\varepsilon^{(r)}, \quad \text{div} \mathbf{z}_\varepsilon^{(r)} = 0, \\ \text{rot} (\eta^\varepsilon)^{-1} \mathbf{w}_\varepsilon^{(r)} + \mathbf{z}_\varepsilon^{(r)} &= i\mathbf{r}, \quad \text{div} \mathbf{w}_\varepsilon^{(r)} = 0. \end{aligned} \right\} \quad (4)$$

Accordingly,

$$\mathbf{u}_\varepsilon^{(r)} = (\eta^\varepsilon)^{-1} \mathbf{w}_\varepsilon^{(r)}, \quad \mathbf{v}_\varepsilon^{(r)} = (\mu^\varepsilon)^{-1} \mathbf{z}_\varepsilon^{(r)}. \quad (5)$$

The results are formulated in terms of the "homogenized" Maxwell system and the "correction" Maxwell system. Let μ^0 be the "effective" matrix (e. g., see [1,2]) corresponding to the elliptic operator $-\text{div} \mu(\mathbf{x})\nabla$. Recall the definition of the (constant positive) matrix μ^0 . Let $\mathbf{e}_1, \mathbf{e}_2, \mathbf{e}_3$ be the standard orthonormal basis in \mathbb{R}^3 , and let $\Phi_j \in H^1_{\text{loc}}(\mathbb{R}^3)$, $j = 1, 2, 3$, be a Γ -periodic solution of the equation $\text{div} \mu(\mathbf{x})(\nabla\Phi_j(\mathbf{x}) + \mathbf{e}_j) = 0$. By $\tilde{\mu}(\mathbf{x})$ we denote the matrix with columns $\mu(\mathbf{x})(\nabla\Phi_j(\mathbf{x}) + \mathbf{e}_j)$, $j = 1, 2, 3$. Then $\mu^0 = |\Omega|^{-1} \int_\Omega \tilde{\mu}(\mathbf{x}) d\mathbf{x}$. The effective matrix η^0 corresponding to the operator $-\text{div} \eta(\mathbf{x})\nabla$ is defined in a similar way. We put $\mathcal{M}^0 = \mathcal{M}(\eta^0, \mu^0)$. Let $(\mathbf{w}_0^{(r)}, \mathbf{z}_0^{(r)})$ be the solution of the "homogenized" Maxwell system

$$(\mathcal{M}^0 - iI) \begin{pmatrix} \mathbf{w}_0^{(r)} \\ \mathbf{z}_0^{(r)} \end{pmatrix} = \begin{pmatrix} 0 \\ \mathbf{r} \end{pmatrix}. \quad (6)$$

We put

$$\mathbf{u}_0^{(r)} = (\eta^0)^{-1} \mathbf{w}_0^{(r)}, \quad \mathbf{v}_0^{(r)} = (\mu^0)^{-1} \mathbf{z}_0^{(r)}. \quad (7)$$

Now we describe the "correction" Maxwell system. Let $F(\mathbf{x})$ be the matrix with columns $\nabla\Phi_j(\mathbf{x})$, $j = 1, 2, 3$. Note that $F(\mathbf{x})$ is a Γ -periodic matrix-valued function with zero mean value. We denote $F^\varepsilon(\mathbf{x}) = F(\varepsilon^{-1}\mathbf{x})$. Let \mathcal{P}_0 be the orthogonal projection in $\mathfrak{G}((\mu^0)^{-1})$ onto J . We put $\mathbf{r}_\varepsilon = \mathcal{P}_0(F^\varepsilon)^*\mathbf{r}$. Then $\mathbf{r}_\varepsilon \in H^{-1}(\mathbb{R}^3; \mathbb{C}^3)$ and $\text{div} \mathbf{r}_\varepsilon = 0$. Let $(\tilde{\mathbf{w}}_\varepsilon^{(r)}, \tilde{\mathbf{z}}_\varepsilon^{(r)})$ be the solution of the "correction" Maxwell system

$$(\mathcal{M}^0 - iI) \begin{pmatrix} \tilde{\mathbf{w}}_\varepsilon^{(r)} \\ \tilde{\mathbf{z}}_\varepsilon^{(r)} \end{pmatrix} = \begin{pmatrix} 0 \\ \mathbf{r}_\varepsilon \end{pmatrix}. \quad (8)$$

We put

$$\tilde{\mathbf{v}}_\varepsilon^{(r)} = (\mu^0)^{-1} \tilde{\mathbf{z}}_\varepsilon^{(r)}. \quad (9)$$

Note that the fields $\tilde{\mathbf{w}}_\varepsilon^{(r)}$, $\tilde{\mathbf{z}}_\varepsilon^{(r)}$, $\tilde{\mathbf{v}}_\varepsilon^{(r)}$ weakly tend to zero in \mathfrak{G} . The reason is that the right-hand side \mathbf{r}_ε in (8) contains the factor F^ε , which weakly tends to zero in $L_{2,\text{loc}}(\mathbb{R}^3)$ (by the "mean value property").

Our main result (as applied to the case $\mathbf{q} = 0$) is the following theorem.

Theorem. *Suppose that Γ -periodic matrix-valued functions $\eta(\mathbf{x})$, $\mu(\mathbf{x})$ satisfy conditions (1). Let $(\mathbf{w}_\varepsilon^{(r)}, \mathbf{z}_\varepsilon^{(r)})$ be the solution of system (4) with $\mathbf{r} \in J$, and let $\mathbf{u}_\varepsilon^{(r)}, \mathbf{v}_\varepsilon^{(r)}$ be defined by (5). Suppose that $(\mathbf{w}_0^{(r)}, \mathbf{z}_0^{(r)})$ is the solution of system (6),*

and let $\mathbf{u}_0^{(r)}, \mathbf{v}_0^{(r)}$ be defined by (7). Suppose that $(\tilde{\mathbf{w}}_\varepsilon^{(r)}, \tilde{\mathbf{z}}_\varepsilon^{(r)})$ is the solution of system (8), and let $\tilde{\mathbf{v}}_\varepsilon^{(r)}$ be defined by (9). Then the following assertions hold.

1°. For the magnetic intensity $\mathbf{v}_\varepsilon^{(r)}$ we have the approximation

$$\|\mathbf{v}_\varepsilon^{(r)} - (\mathbf{1} + F^\varepsilon)(\mathbf{v}_0^{(r)} + \tilde{\mathbf{v}}_\varepsilon^{(r)})\|_{\mathfrak{G}} \leq C\varepsilon\|\mathbf{r}\|_{\mathfrak{G}}, \quad 0 < \varepsilon \leq 1. \quad (10)$$

2°. As $\varepsilon \rightarrow 0$, $\mathbf{v}_\varepsilon^{(r)}$ weakly tends in \mathfrak{G} to $\mathbf{v}_0^{(r)}$, and $\text{rot } \mathbf{v}_\varepsilon^{(r)}$ weakly tends in \mathfrak{G} to $\text{rot } \mathbf{v}_0^{(r)}$.

3°. For the magnetic induction vector $\mathbf{z}_\varepsilon^{(r)}$ we have the approximation

$$\|\mathbf{z}_\varepsilon^{(r)} - (\mathbf{1} + G^\varepsilon)(\mathbf{z}_0^{(r)} + \tilde{\mathbf{z}}_\varepsilon^{(r)})\|_{\mathfrak{G}} \leq C\varepsilon\|\mathbf{r}\|_{\mathfrak{G}}, \quad 0 < \varepsilon \leq 1, \quad (11)$$

where $G(\mathbf{x}) := \tilde{\mu}(\mathbf{x})(\mu^0)^{-1} - \mathbf{1}$ is a Γ -periodic matrix-valued function with zero mean value, and $G^\varepsilon(\mathbf{x}) = G(\varepsilon^{-1}\mathbf{x})$.

4°. As $\varepsilon \rightarrow 0$, $\mathbf{z}_\varepsilon^{(r)}$ weakly tends in \mathfrak{G} to $\mathbf{z}_0^{(r)}$.

5°. As $\varepsilon \rightarrow 0$, the electric field intensity $\mathbf{u}_\varepsilon^{(r)}$ weakly tends in \mathfrak{G} to $\mathbf{u}_0^{(r)}$. For $\text{rot } \mathbf{u}_\varepsilon^{(r)} = \mathbf{i}\mathbf{r} - \mathbf{z}_\varepsilon^{(r)}$ we have approximation in the \mathfrak{G} -norm (see (11)).

6°. As $\varepsilon \rightarrow 0$, the electric displacement vector $\mathbf{w}_\varepsilon^{(r)}$ weakly tends in \mathfrak{G} to $\mathbf{w}_0^{(r)}$.

The case where $\mathbf{r} = 0$ can be treated in a similar way. For the fields $\mathbf{u}_\varepsilon^{(q)}$ and $\mathbf{w}_\varepsilon^{(q)}$ we obtain approximations in the \mathfrak{G} -norm similar to (10), (11), while for $\mathbf{v}_\varepsilon^{(q)}$, $\mathbf{z}_\varepsilon^{(q)}$ we have only the weak convergence.

In the case where permeability is constant: $\mu = \mu^0$, the results simplify. In this case the solutions of the "correction" system (8) are trivial, and we have

$$\|\mathbf{v}_\varepsilon^{(r)} - \mathbf{v}_0^{(r)}\|_{\mathfrak{G}} \leq C\varepsilon\|\mathbf{r}\|_{\mathfrak{G}}, \quad \|\mathbf{z}_\varepsilon^{(r)} - \mathbf{z}_0^{(r)}\|_{\mathfrak{G}} \leq C\varepsilon\|\mathbf{r}\|_{\mathfrak{G}}, \quad 0 < \varepsilon \leq 1. \quad (12)$$

For $\mu(\mathbf{x}) = \mathbf{1}$ this result has been obtained before in [4]. If $\eta = \eta^0$, then for $\mathbf{u}_\varepsilon^{(q)}$, $\mathbf{w}_\varepsilon^{(q)}$ we have similar results.

REFERENCES

- [1] A. Bensoussan A., J.-L. Lions, G. Papanicolaou, *Asymptotic analysis for periodic structures*, Stud. Math. Appl., vol. 5, North-Holland Publishing Co., Amsterdam-New York, 1978.
- [2] V. V. Zhikov, S. M. Kozlov and O. A. Oleinik, *Homogenization of differential operators*, "Nauka", Moscow, 1993; English transl., Springer-Verlag, Berlin, 1994.
- [3] M. Sh. Birman and T. A. Suslina, *Threshold effects near the lower edge of the spectrum for periodic differential operators of mathematical physics*, Systems, Approximations, Singular Integral Operators and Related Topics (Bordeaux, 2000), Oper. Theory Adv. Appl., vol. 129, Birkhäuser, Basel, 2001, pp. 71–107.
- [4] M. Sh. Birman and T. A. Suslina, *Second order periodic differential operators. Threshold properties and homogenization*, Algebra i Analiz **15** (2003), no. 5, 1–108; English transl., St. Petersburg Math. J. **15** (2004), no. 5, 1–77.
- [5] T. A. Suslina, *On the homogenization of the periodic Maxwell system*, Funct. Anal. Appl. **38** (2004), no. 3, 234–237.
- [6] T. A. Suslina, *Homogenization of a stationary periodic Maxwell system*, Algebra i Analiz **16** (2004), no. 5, 162–244; English transl., St. Petersburg Math. J. **16** (2005), no. 5.

Bound States and Essential Spectrum

DAVID DAMANIK

Given a Schrödinger operator $H_V = -\Delta + V$ in $L^2(\mathbb{R}^d)$ or $h_V = \Delta + V$ in $\ell^2(\mathbb{Z}^d)$, a basic problem is to study the discrete spectrum and the essential spectrum. In recent years, several papers have uncovered a surprising connection between these two parts of the spectrum.

We first consider operators with empty discrete spectrum. The following theorem was shown by Killip and Simon [5]:

Theorem 1. *Suppose $\mathcal{H} = \ell^2(\mathbb{Z})$. Then $\sigma(h_V) \subseteq [-2, 2]$ implies $V \equiv 0$.*

In particular, every potential that does not vanish identically must produce spectrum outside of the free spectrum, $[-2, 2]$. Note that this result holds without any a priori assumption on V . The proof of Theorem 1 given in [5] relies on sum rules and is, to some extent, a by-product of their general study culminating in a characterization of all (half-line) Jacobi matrices that are Hilbert-Schmidt perturbations of the free operator in terms of properties of the spectral measure.

A more direct and elementary proof based on suitable choices of test functions was given in [1]. Moreover, it was possible to extend the result to two dimensions:

Theorem 2. *Suppose $\mathcal{H} = \ell^2(\mathbb{Z}^d)$ with $d = 1$ or 2 . Then $\sigma(h_V) \subseteq [-2d, 2d]$ implies $V \equiv 0$.*

It was also shown in [1] that $\sigma_{\text{ess}}(h_V) \subseteq [-2d, 2d]$ implies $V \rightarrow 0$. Both statements fail in dimensions $d \geq 3$.

On the half-line, the following example was discussed in [1]. Consider the operator h_V in $\ell^2(\mathbb{Z}_+)$ with potential $V(n) = (-1)^n/n$. Then, h_V has spectrum $[-2, 2]$. This shows that not only can $V \equiv 0$ fail under the assumption $\sigma(h_V) \subseteq [-2, 2]$ for half-line operators, it is also not immediately clear what one can say about the spectral type inside $[-2, 2]$.

Half-line operators were studied in [2], where the following theorem was proven:

Theorem 3. *Suppose $\mathcal{H} = \ell^2(\mathbb{Z}_+)$. Then $\sigma(h_V) \subseteq [-2, 2]$ implies $\sigma_{\text{sing}}(h_V) = \emptyset$.*

The main steps in the proof of Theorem 3 are as follows: First map the spectral measure to the unit circle via $E = z + z^{-1}$ and find relations between the potential and the Verblunsky coefficients of the associated measure. Then use this connection to find bounds on the potential. Finally, use these bounds to show that there cannot be any embedded singular spectrum. For example, it is shown that under the assumption $\sigma(h_V) \subseteq [-2, 2]$, V may be written as

$$V(n) = W(n) - W(n-1) + Q(n),$$

where

$$Q \in \ell^1 \text{ and } \sum_{n=1}^N nW(n)^2 \leq \frac{1}{4} \log N + C.$$

The continuum case is also studied in [2]. Note that the unitary $[U\phi](n) = (-1)^n\phi(n)$ conjugates h_{-V} and $-h_V$. Therefore, $\sigma(h_V) \subseteq [-2, 2]$ is equivalent to

the two conditions $\sigma(h_{\pm V}) \subseteq [-2, \infty)$. Thus, the following theorem from [2] is the natural continuum analogue of Theorem 3:

Theorem 4. *Suppose $\mathcal{H} = L^2(\mathbb{R}_+)$ and $V \in \ell^\infty(L^2)$. Then $\sigma(H_{\pm V}) \subseteq [0, \infty)$ implies $\sigma_{\text{sing}}(H_V) = \emptyset$.*

Given this observation, it is natural to ask whether the $\ell^2(\mathbb{Z})$ and $\ell^2(\mathbb{Z}^2)$ results have continuum analogues. This is indeed the case, as shown in [3]:

Theorem 5. *Let $d = 1$ or 2 . Suppose that $Q \in L^2_{\text{loc}}(\mathbb{R}^d)$ and the operator H_Q has a bounded positive ground state. If $V \in L^2_{\text{loc}}(\mathbb{R}^d)$ and both $H_{Q \pm V}$ are bounded below by the ground state energy of H_Q , then $V \equiv 0$.*

For $Q \equiv 0$, this gives the continuum analogue of Theorem 2 (choose $\psi \equiv 1$ as the bounded positive ground state), but it also applies to periodic Q , for example.

Thus, one has a good understanding of cases without bound states. If a perturbation introduces only finitely many bound states, one may still hope for strong restrictions on V and the spectral type inside the essential spectrum. In fact, Theorems 3 and 4 extend quite easily to the case of finitely many bound states. Thus, on the half-line, finiteness of the number of bound states implies the absence of embedded singular spectrum, as shown in [2]. For operators on the line, the corresponding results were shown in [3]. The problem is open in two dimensions.

The following example shows that an extension to operators with infinitely many bound states could be involved. On the half-line, the operator with Wigner-von Neumann-type potential $V(n) = (1 + \varepsilon)(-1)^n/n + O(1/n^2)$, $\varepsilon > 0$, has an embedded eigenvalue and the discrete eigenvalues decay exponentially; see [3]. In particular, the finiteness of bound state moments is not sufficient to exclude embedded singular spectrum. On the other hand, positive results are obtained in [4]. For example, if the p -th bound state moment is finite, then the embedded singular spectrum must be supported on a set of Hausdorff dimension $4p$.

REFERENCES

- [1] D. Damanik, D. Hundertmark, R. Killip, and B. Simon, *Variational estimates for discrete Schrödinger operators with potentials of indefinite sign*, Commun. Math. Phys. **238** (2003), 545–562.
- [2] D. Damanik and R. Killip, *Half-line Schrödinger operators with no bound states*, to appear in Acta Math.
- [3] D. Damanik, R. Killip and B. Simon, *Schrödinger operators with few bound states*, preprint (arXiv/math-ph/0409074).
- [4] D. Damanik and C. Remling, in preparation.
- [5] R. Killip and B. Simon, *Sum rules for Jacobi matrices and their applications to spectral theory*, Ann. of Math. **158** (2003), 253–321.

Properties of Coulombic wavefunctions and their electron density

MARIA HOFFMANN-OSTENHOF

(joint work with Søren Fournais, Thomas Hoffmann-Ostenhof and Thomas Østergaard Sørensen)

Let H be the non-relativistic Schrödinger operator of an N -electron atom with nuclear charge Z and the nucleus fixed in the origin, given by

$$(1) \quad H = -\Delta + V = \sum_{j=1}^N \left(-\Delta_j - \frac{Z}{|x_j|} \right) + \sum_{1 \leq i < j \leq N} \frac{1}{|x_i - x_j|}$$

where $x_j \in \mathbb{R}^3$, $1 \leq j \leq N$, are the electron coordinates and the Δ_j the associated Laplacians. Let $\psi \in L^2(\mathbb{R}^{3N})$ be an eigenfunction of H with eigenvalue E . By a classical result of Kato [5], ψ is locally Lipschitz and real analytic away from the singularities of the potential V .

The main result in [3] is a representation result for electronic wavefunctions of atoms and molecules which is stated here for simplicity for the atomic case:

Theorem

Let ψ be as above and let

$$\mathcal{F} = e^{F_2 + F_3}$$

with

$$F_2 = -\frac{Z}{2} \sum_{i=1}^N |x_i| + \frac{1}{4} \sum_{1 \leq i < j \leq N} |x_i - x_j|$$

and

$$F_3 = c_0 Z \sum_{1 \leq i < j \leq N} \langle x_i | x_j \rangle \ln(|x_i|^2 + |x_j|^2), \quad c_0 = \frac{2 - \pi}{12\pi}.$$

Then

$$\psi = \mathcal{F}\Phi \text{ with } \Phi \in C^{1,1}(\mathbb{R}^{3N}).$$

From this and earlier results of the present authors certain properties of ψ and the associated 1-electron density

$$\rho(x) = \int_{\mathbb{R}^{3N-3}} |\psi|^2(x, x') dx', \quad x \in \mathbb{R}^3$$

can be shown (work in progress): In 1957 Kato, [5], analyzed the behaviour of ψ in an averaged sense near two particle coalescence points, (Kato's cusp conditions). Generalizations of such "cusp properties" are investigated. By a cusp condition (resp. property) we will understand a condition a solution ψ has to satisfy near a point, where the potential in (1) is singular. In an L^∞ -sense such properties concerning second order partial derivatives of ψ are given in [3]. In progress is work on cusp conditions for the 1-electron density ρ . It can be shown that for some $\vec{c} \in \mathbb{R}^3$

$$\left. \frac{\partial \rho(r\omega)}{\partial r} \right|_{r=0} = \lim_{r \downarrow 0} \frac{\partial}{\partial r} \rho(r\omega) = -Z\rho(0) + \langle \omega | \vec{c} \rangle$$

where $x = r\omega$, so that $\omega \in \mathbb{S}^2$.

Other investigations concern the regularity properties of ρ , respectively, of the spherically averaged density $\tilde{\rho} = \tilde{\rho}(r)$. It is known [1, 2] that ρ is smooth and even real analytic away from the origin. An open question is the regularity of $\tilde{\rho}$ near the origin \mathcal{O} . So far only $C^2([0, \infty))$ was known, [4], and this can be extended, via appropriate estimates derived in [3], to C^3 .

Another interesting question is whether ρ is strictly positive in \mathbb{R}^3 . Of course this cannot be true in general, since it fails for excited states of the Hydrogen atom. It is well known that the mathematical groundstate ψ satisfies $|\psi| > 0$ in \mathbb{R}^{3N} and therefore the associated density ρ is strictly positive.

We investigate the spherically averaged density associated to a groundstate of an atom in some symmetry subspace and are going to show that it is strictly positive and we will also give an explicit lower bound to $\rho(0)$. (Note that for these considerations we use the symmetrized (physical) density ρ instead of the one defined above.)

Whether $\tilde{\rho}(r)$ is monotonically decreasing is an open problem for decades. This monotonicity is expected to hold for groundstate densities, but not known even for the bosonic case like the Helium groundstate in spite of overwhelming numerical evidence. So far it is only known in a sufficiently small neighborhood of the origin and sufficiently far away from the nucleus.

REFERENCES

- [1] Søren Fournais, Maria Hoffmann-Ostenhof, Thomas Hoffmann-Ostenhof, Thomas Østergaard Sørensen. The electron density is smooth away from the nuclei. *Commun. Math. Phys.*, 228, 401–415, 2002.
- [2] ———. Analyticity of the density of electronic wavefunctions. *Ark. Mat.*, 42, 87–106, 2004.
- [3] ———. Sharp regularity results for Coulombic many-electron wave functions. *Commun. Math. Phys.* in press, 2004.
- [4] Maria Hoffmann-Ostenhof, Thomas Hoffmann-Ostenhof, Thomas Østergaard Sørensen. Electron wave functions and densities for atoms. *Ann. Henri Poincaré* 2, 77–100, 2001.
- [5] Tosio Kato. On the eigenfunctions of many-particle systems in quantum mechanics. *Commun. Pure and Appl. Math.* 10, 151–177, 1957.

On a Magnetic Hardy Inequality in The Waveguide

HYNEK KOVAŘÍK

(joint work with Denis Borisov and Tomas Ekholm)

It is well known that the classical Hardy inequality fails to hold in dimensions one and two. This is closely related to the fact, that arbitrarily small attractive potential perturbation produces at least one negative eigenvalue of the Schrödinger operator on $L^2(\mathbb{R}^d)$, $d = 1, 2$. As a consequence, the threshold of the spectrum of the Laplacian in the so called quantum waveguide, i.e. in a two-dimensional strip Ω with Dirichlet boundary conditions, is unstable under any local enlargement or bending of the waveguide, see [BGRS], [EŠ], [GJ].

On the other hand, in 1999 Laptev and Weidl proved a modified version of the Hardy inequality in \mathbb{R}^2 for the quadratic form of a magnetic Schrödinger operator

$$(1) \quad \text{Const} \int_{\mathbb{R}^2} \frac{|u(x)|^2}{1+|x|^2} dx \leq \int_{\mathbb{R}^2} |(-i\nabla + A)u(x)|^2 dx,$$

and gave a sharp result for the case of Aharonov-Bohm field. See [LW] for details. This work was later extended in [B] to multiple Aharonov-Bohm magnetic potentials, see also [EL]. Recently another generalization of the result by Laptev and Weidl was obtained in [BLS].

In our model we study the spectrum of the magnetic Schrödinger operator $(-i\nabla + A)^2$ on $L^2(\Omega)$ with $\Omega = \mathbb{R} \times (0, \pi)$. Essential difference to the case treated in [LW] is that due to the Dirichlet boundary conditions the spectrum starts from 1. Consequently inequality (1) becomes trivial. Therefore we shall subtract the threshold of the spectrum and prove a Hardy-inequality in the form

$$(2) \quad \text{Const} \int_{\mathbb{R} \times (0, \pi)} \frac{|u(x)|^2}{1+x_1^2} dx \leq \int_{\mathbb{R} \times (0, \pi)} (|(-i\nabla + A)u(x)|^2 - |u(x)|^2) dx,$$

for all u in the magnetic Sobolev space $H_{0,A}^1(\mathbb{R} \times (0, \pi))$. We prove this inequality for the magnetic Schrödinger operator with a locally bounded field.

As an application of inequality (2) we show that the threshold of the spectrum of the corresponding magnetic Schrödinger operator is stable under local geometrical perturbations of the waveguide as well as under local perturbations of the boundary conditions. In the first case it is shown that a sufficiently weak enlargement of the waveguide, depending on the magnetic field, will not produce any discrete spectrum of the operator $(-i\nabla + A)^2$. In a similar way we note that the discrete spectrum of $(-i\nabla + A)^2$ in a mildly curved waveguide stays empty as long as the corresponding curvature and its first derivative are small enough.

In the second model we consider a situation where the Dirichlet boundary condition is switched to magnetic Neumann on a fixed segment of the length $2l$ of the boundary of Ω . Such a perturbation is stronger than the geometrical perturbations of the boundary mentioned above and therefore a different approach is needed in order to establish the desired stability result. Using a similar integral inequality to (2) we are able to show that it suffices to prove the non-existence of discrete eigenvalues of the one-dimensional Schrödinger operator

$$A = -\frac{d^2}{dx^2} + V,$$

where V is a sum of the purely attractive potential well of the width $2l$ and a small, but fixed positive potential. We conclude that the discrete spectrum of A and consequently also the discrete spectrum of the corresponding magnetic Schrödinger operator stays empty provided l is small enough.

This talk has been based on two papers, [EK] and [BEK], obtained in the collaboration with T.Ekholm and D.Borisov, T.Ekholm respectively.

REFERENCES

- [B] A.A. Balinsky: Hardy type inequalities for Aharonov-Bohm magnetic potentials with multiple singularities, *mp-arc 02-416*.
- [BLS] A.A. Balinsky, A. Laptev and A.V. Sobolev: Generalized Hardy inequality for the magnetic Dirichlet forms, *J. Stat. Phys.*, **116**, no. 1/4, (2004), 507–521.
- [BEK] D. Borisov, T. Ekholm and H. Kovařík: Spectrum of the magnetic Schrödinger operator in a waveguide with combined boundary conditions, to appear in *Ann. Henri Poincaré*. Preprint: math-ph/0405034.
- [BGRS] W. Bulla, F. Gesztesy, W. Renger and B. Simon: Weakly coupled bound states in quantum waveguides, *Proc. Amer. Math. Soc.* **125** (1997), no. 5, 1487–1495.
- [EK] T. Ekholm and H. Kovařík, Stability of the magnetic Schrödinger operator in a waveguide, to appear in *Comm. in Part. Diff. Eq.* Preprint: math-ph/0404069.
- [EŠ] P. Exner and P. Šeba: Bound states in curved quantum waveguides, *J. Math. Phys.* **30** (1989), 2574–2580.
- [EL] W.D. Evans and R.T. Lewis: On the Rellich inequality with magnetic potentials, *mp-arc 04-93*.
- [GJ] J. Goldstone and R.L. Jaffe: Bound states in twisting tubes, *Phys. Rev.* **B45** (1992), 14100–14107.
- [LW] A. Laptev and T. Weidl: Hardy inequalities for magnetic Dirichlet forms, *Oper. Theory Adv. Appl.* **108** (1999) 299–305.

Spectral asymptotics for the Landau Hamiltonian and logarithmic capacity

ALEXANDER PUSHNITSKI

(joint work with Nikolai Filonov)

Consider the operator

$$H_0 = \left(-i \frac{\partial}{\partial x} + \frac{B}{2} y \right)^2 + \left(-i \frac{\partial}{\partial y} - \frac{B}{2} x \right)^2 \quad \text{in } L^2(\mathbb{R}^2, dx dy),$$

where $B > 0$ is the strength of the constant magnetic field. The spectrum of H_0 consists of the eigenvalues (known as Landau levels) $\Lambda_q = B(2q+1)$, $q = 0, 1, 2, \dots$; each of these eigenvalues has infinite multiplicity.

Let $V \geq 0$ be a perturbation potential such that $V \in L^\infty(\mathbb{R}^2)$, and $\Omega = \text{supp}(V)$ is compact. We consider the spectrum of the operators $H_\pm = H_0 \pm V$. It is well known that $\sigma_{\text{ess}}(H_\pm) = \sigma_{\text{ess}}(H_0) = \cup_{q=0}^\infty \Lambda_q$. Moreover, for any q the eigenvalues of H_+ can accumulate to Λ_q only from the right, and eigenvalues of H_- can accumulate to Λ_q only from the left.

Let us enumerate the eigenvalues of H_- in $(-\infty, \Lambda_0)$ in the ascending order (counting multiplicities):

$$\lambda_1^- \leq \lambda_2^- \leq \dots \leq \lambda_n^- \leq \dots < \Lambda_0.$$

Similarly, let us enumerate the eigenvalues of H_+ in (Λ_0, Λ_1) in the descending order (counting multiplicities):

$$\Lambda_0 < \dots \leq \lambda_n^+ \leq \dots \leq \lambda_2^+ \leq \lambda_1^+ < \Lambda_1.$$

For a bounded Borel set $C \subset \mathbb{R}^2$, we denote by $\text{Cap } C$ the logarithmic capacity of C (see [2]). Denote

$$\rho(V) = \text{Cap } \Omega, \quad \Omega = \text{supp } V,$$

$$\rho_-(V) = \inf\{\text{Cap } C \mid C \subset \mathbb{R}^2 \text{ is a bounded Borel set, } \int_{\mathbb{R}^2 \setminus C} V(x, y) dx dy = 0\}.$$

Clearly, $\rho_-(V) \leq \rho(V)$.

Theorem *Assume that $\rho_-(V) = \rho(V)$. Then one has the asymptotics:*

$$(*) \quad \lambda_n^\pm - \Lambda_0 = \pm \frac{1}{n!} \left(\frac{B}{2} \rho(V)^2 \right)^{n+o(n)}, \quad n \rightarrow \infty.$$

Remarks

(1) The asymptotics (*) should be understood as

$$\log(\pm(\lambda_n^\pm - \Lambda_0)n!) = n \log \left(\frac{B}{2} \rho(V)^2 \right) + o(n), \quad n \rightarrow \infty.$$

(2) An elementary calculation shows that the asymptotics (*) is equivalent to

$$\lambda_{n+k}^\pm - \Lambda_0 = \pm \frac{1}{n!} \left(\frac{B}{2} \rho(V)^2 \right)^{n+o(n)}, \quad n \rightarrow \infty,$$

for any integer k .

(3) We also have a way of treating the eigenvalue asymptotics near higher Landau levels Λ_q , $q \geq 1$, but at present this construction requires more restrictive assumptions on V and slightly more technical arguments, so we do not include it in this preliminary report.

For $t > 0$, let us define $N_-(t)$ as the total number of eigenvalues of H_- (counting multiplicities) in the interval $(-\infty, \Lambda_0 - t)$. Similarly, for $0 < t < 2B$, let us define $N_+(t)$ as the total number of eigenvalues of H_+ in $(\Lambda_0 + t, \Lambda_1)$. An elementary calculation shows that (*) is equivalent to

$$N_\pm(t) = \frac{|\log t|}{(\log|\log t|)^2} \left(\log|\log t| + \log \log|\log t| + \log \left(\frac{B}{2} \rho(V)^2 \right) + 1 + o(1) \right),$$

as $t \rightarrow +0$. In the papers [4] and [3], the asymptotics

$$N_\pm(t) = \frac{|\log t|}{(\log|\log t|)}(1 + o(1)), \quad t \rightarrow +0$$

was obtained.

The proof is based on the following ideas. First, as in the papers [4] and [3], we reduce the question to the asymptotics of the eigenvalues of the auxiliary operator P_0VP_0 . Here P_0 is the spectral projection of H_0 , corresponding to the first Landau level Λ_0 . Next, the eigenvalues of P_0VP_0 are identified with the singular numbers of the embedding $F \subset L^2(\mathbb{R}^2, \tilde{V}(x, y) dx dy)$, where F is the so-called Fock class, and $\tilde{V}(x, y) = V(x, y)e^{-x^2-y^2}$. Using the techniques of [6], the singular numbers of this embedding are then expressed in terms of the asymptotics for a certain sequence

of orthogonal polynomials. Finally, application of the “regularity criteria” of [5] yields the required result.

REFERENCES

- [1] E. Hille, *Analytic function theory. Vol. II*, Ginn and Co., Boston, Mass., 1962
- [2] N. S. Landkof, *Foundations of modern potential theory*, Springer, New York, 1972
- [3] M. Melgaard and G. Rozenblum, *Eigenvalue asymptotics for weakly perturbed Dirac and Schrödinger operators with constant magnetic fields of full rank*, *Comm. Partial Differential Equations* **28** (2003), no. 3–4, 697–736.
- [4] G. D. Raikov and S. Warzel, *Quasi-classical versus non-classical spectral asymptotics for magnetic Schrödinger operators with decreasing electric potentials*, *Rev. Math. Phys.* **14** (2002), no. 10, 1051–1072
- [5] H. Stahl and V. Totik, *General orthogonal polynomials*, Cambridge Univ. Press, Cambridge, 1992.
- [6] O. G. Parfënov, *The widths of some classes of entire functions*, *Mat. Sb.* **190** (1999), no. 4, 87–94; translation in *Sb. Math.* **190** (1999), no. 3–4, 561–568.

Gaussian extremizers for the Strichartz inequality in one and two dimensions

DIRK HUNDERTMARK

(joint work with Vadim Zharnitsky)

We show that in dimension one and two the only maximizers for the homogenous Strichartz inequality for the free Schrödinger evolution are Gaussians.

More precisely, let u be the solution to the free Schrödinger equation

$$(1) \quad i\partial_t u = \Delta u$$

with initial condition $u(0) = f \in L^2(\mathbb{R}^2)$. It is, of course, given by

$$(2) \quad u(t, x) = (e^{-it\Delta} f)(x)$$

where $e^{-it\Delta}$ is defined, for example, by the functional calculus. Since, for fixed time t , $e^{-it\Delta}$ is a unitary operator on $L^2(\mathbb{R}^d)$, one immediately sees that $u \in L_t^\infty(L^2(\mathbb{R}^d))$. But, in fact, due to the dispersive nature of the free Schrödinger equation, the solution u , as a function of space-time, obeys the stronger L^p -bound

$$(3) \quad \|u\|_{L^p(\mathbb{R} \times \mathbb{R}^d)} \leq S_d \|f\|_{L^2(\mathbb{R}^d)}$$

where $p = p(d) = 2 + \frac{4}{d}$. This was first shown by Strichartz [6] who followed the L^p restriction proof of Stein-Tomas. Later simplified proofs were given by Ginibre and Velo [4], see also [2, 7].

The sharp value of S_d , i.e., the quantity

$$(4) \quad S_d = \sup_{f \neq 0} \frac{\|u\|_{L^p(\mathbb{R} \times \mathbb{R}^d)}}{\|f\|_{L^2(\mathbb{R}^d)}}$$

has been unknown until very recently. In fact, even the existence of maximizers for (4), that is, functions $f_* \neq 0$ such one has equality in (4),

$$(5) \quad S_d = \frac{\|e^{-it\Delta} f_*\|_{L^p(\mathbb{R}^{d+1})}}{\|f_*\|_{L^2(\mathbb{R})}},$$

has been only recently established. By using an elaborate application of Lions' concentration compactness method, Markus Kunze showed in [5] that (4) has a maximizer in one dimension. His proof does not, however, provide any explicit information about the maximizer nor the value of S_1 . The reason why even the existence of maximizers has not been known until recently is the invariance of the Strichartz inequality under the rather large group of Galilei transformations and scaling. This makes the usual existence proof for maximizers via minimizing sequences very hard, since they can very easily converge weakly to zero. The usual method to circumvent this is the concentration compactness principle, however, in this setting it has to be used twice, first in Fourier space, then in real space.

Very recently, Damiano Foschi [3] gave a proof of the Strichartz inequality in one and two dimensions, which yields the sharp constant. He showed

Theorem 1 (Foschi 2004, [3]). *The sharp constants for the Strichartz inequality in one and two dimensions are $S_1 = 12^{-1/12}$ and $S_2 = 2^{-1/2}$, respectively. Moreover, if the initial condition f is given by a Gaussian, then one has equality in the Strichartz inequality.*

However, the existence of non-Gaussian maximizers was not ruled out in [3]. The main purpose of this note is to give a simple argument which shows that at least in one and two dimensions the *only maximizers* in the Strichartz inequality are Gaussians. More precisely, we have the following

Theorem 2 (Gaussian maximizers). *Let $d = 1$ or 2 . The function $f_* \in L^2(\mathbb{R}^d)$ is a maximizer for the Strichartz inequality (3), that is, (4) holds, if and only if f_* is a Gaussian. More precisely, there exists $A \in \mathbb{C}$, $\lambda > 0$, $\mu \in \mathbb{R}$, $a \in \mathbb{R}^d$, and $b \in \mathbb{C}^d$ such that*

$$(6) \quad f_*(x) = A e^{(-\lambda + i\mu)|x-a|^2 + b \cdot x}.$$

The key for our proof is the following representation theorem. It shows that the Strichartz estimate follows from a simple bound on a *linear* operator and, moreover, gives a geometric criterion for the maximizer in the Strichartz inequality. For $f \in L^2(\mathbb{R}^d)$, denote by $f \otimes f$ be the usual tensor product, $\mathbb{R}^d \times \mathbb{R}^d \ni x = (x_1, x_2) \rightarrow f \otimes f(x) := f(x_1)f(x_2)$. Similarly for the triple tensor product $f \otimes f \otimes f$. Furthermore, let $P_1 : L^2(\mathbb{R}^3) \rightarrow L^2(\mathbb{R}^3)$ be the orthogonal projection operator onto the subspace consisting of functions $F \in L^2(\mathbb{R}^3)$ which are symmetric under rotations of \mathbb{R}^3 keeping the $(1, 1, 1)$ direction fixed. And similarly, let $P_2 : L^2(\mathbb{R}^4) \rightarrow L^2(\mathbb{R}^4)$ be the orthogonal projection operator onto functions $F \in L^2(\mathbb{R}^3)$ which are symmetric under rotations of \mathbb{R}^4 fixing both the $(1, 0, 1, 0)$ and $(0, 1, 0, 1)$ direction. With this, we have

Theorem 3. *Let $f \in L^2(\mathbb{R}^d)$.*

a) In dimension one,

$$\int_{\mathbb{R}} \int_{\mathbb{R}} |e^{-it\Delta} f(x)|^6 dx dt = \frac{1}{2\sqrt{3}} \langle \hat{f} \otimes \hat{f} \otimes \hat{f}, P_1(\hat{f} \otimes \hat{f} \otimes \hat{f}) \rangle_{L^2(\mathbb{R}^3)}$$

b) In dimension two,

$$\int_{\mathbb{R}} \int_{\mathbb{R}^2} |e^{-it\Delta} f(x)|^4 dx dt = \frac{1}{4} \langle \hat{f} \otimes \hat{f}, P_2(\hat{f} \otimes \hat{f}) \rangle_{L^2(\mathbb{R}^4)},$$

where \hat{f} is the (space) Fourier transform of f .

One immediately gets the sharp Strichartz inequality, using that any projection operator operator is bounded by the identity. One also sees that, in order to have equality in the Strichartz inequality, the function $\hat{f} \otimes \hat{f} \otimes \hat{f}$ must be in the range of P_1 in dimension one, and similarly for the two-dimensional case. In other words, for any one-dimensional maximizer f of the Strichartz inequality, the function $\hat{f} \otimes \hat{f} \otimes \hat{f}$ is invariant under rotations of \mathbb{R}^3 which keep the $(1, 1, 1)$ direction fixed. Similarly, for any two-dimensional maximizer f , the function $\hat{f} \otimes \hat{f}$ is invariant under rotations of \mathbb{R}^4 which keep both the $(1, 0, 1, 0)$ and $(0, 1, 0, 1)$ directions fixed. This is obviously the case if \hat{f} , and hence f , is a Gaussian, and a simple proof, mimicked after a result by Carlen [1], shows that this geometric condition forces f to be a Gaussian.

REFERENCES

- [1] E. Carlen, *subadditivity of Fisher's information and logarithmic Sobolev inequalities*, J. Funct. Anal. **101**, no. 1, 194–211.
- [2] Thierry Cazenave, *Semilinear Schrödinger equations*, Courant Lecture Notes in Mathematics **10**, AMS, Providence, Rhode Island.
- [3] D. Foschi, *Maximizers for the Strichartz inequality*, Preprint, Mai 2004.
- [4] J. Ginibre and G. Velo, *The global Cauchy problem for the nonlinear Schrödinger equation*, Ann. Inst. H. Poincaré Anal. Non Linéaire **2** (1985), 3009–327.
- [5] M. Kunze, *On the existence of a maximizer for the Strichartz inequality*, Commun Math. Phys. **243** (2003), no. 1, 137–162.
- [6] M. Strichartz, *Restrictions of fourier transforms to quadratic surfaces and decay of solutions of wave equations*, Duke math. J. **44** (1977), 705–714.
- [7] C. Sulem, P.-L. Sulem, *The non-linear Schrödinger equation. Self-focusing and wave collapse*, Applied Mathematical Sciences, 139. Springer-Verlag, New York, 1999.

Some Results on the Spectra of Periodic Landau Operators

DANIEL M. ELTON

The *periodic Landau operator* on \mathbb{R}^d is the magnetic Schrödinger operator

$$(1) \quad H_{B,V} = (D - A)^2 + V,$$

where $D = -i\nabla$, V is a periodic potential (with respect to some lattice Λ) and the magnetic field $B = \nabla \times A$ is constant. In this talk we are interested in spectral problems related to $H_{B,V}$ when $d = 2, 3$. For technical convenience we take $B =$

$(0, 0, \beta)$ with $\beta > 0$ when $d = 3$; when $d = 2$ this reduces to $B = \beta$. We will also assume that the lattice of periods is $\Lambda = (2\pi\mathbb{Z})^d$.

1. REDUCTION OF THE OPERATOR

To study the spectrum $\sigma(H_{B,V})$ we firstly use a metaplectic transformation to replace (1) with a unitarily equivalent operator; when $d = 2$ we get

$$(1) \quad D_x^2 + (\beta x)^2 + \text{Op}^w(V(x + y/\beta, -\eta - \xi/\beta))$$

acting in $L^2(\mathbb{R}^2)$, where $\text{Op}^w(p)$ denotes the Weyl-quantised pseudo-differential operator (on \mathbb{R}^2) with symbol $p(x, y, \xi, \eta)$. In general (1) is a harmonic oscillator with a free variable, perturbed by an oscillatory pseudo-differential operator.

Although V and B are periodic functions, $H_{B,V}$ is not a periodic operator (owing to the presence of the magnetic potential A in (1)). The Bloch (or Floquet) techniques commonly used for periodic spectral problems are not directly available to study the spectrum $\sigma(H_{B,V})$. A partial remedy is available via a symmetry group consisting of “magnetic Bloch-transformations”; however this is only useful under the *flux rationality assumption*:

$$(FR) \quad \beta = |B| = \frac{p}{2\pi q} \quad \text{for some } p, q \in \mathbb{N}.$$

Under this condition, $H_{B,V}$ is unitarily equivalent to a direct integral

$$(2) \quad \int_{[0,1)}^{\oplus} dk_1 \int_{S^1}^{\oplus} dk_2 \mathcal{H}(k_1, k_2)$$

with fibre operator $\mathcal{H}(k_1, k_2) = (D_x^2 + (\beta x)^2) \otimes \mathcal{I}_p + \mathcal{A}(k_1, k_2)$ acting on $\bigoplus_{j=0}^{p-1} L^2(\mathbb{R})$; the potential V has become a $p \times p$ matrix of oscillatory pseudo-differential operators $\mathcal{A}(k_1, k_2)$. Since the spectrum of $(D_x^2 + (\beta x)^2) \otimes \mathcal{I}_p$ consists of discrete eigenvalues (the eigenvalues of the 1-dimensional harmonic oscillator, each with multiplicity p), we obtain a band gap picture for $\sigma(H_{B,V})$ (the bands are simply the ranges of the eigenvalues of $\mathcal{H}(k_1, k_2)$ considered as functions of the parameters k_1, k_2).

The above discussion modifies in the obvious way for the case $d = 3$.

2. DIMENSION $d = 3$

When $V \equiv 0$ a straightforward calculation shows $\sigma(H_{B,0}) = [\beta, \infty)$ and this spectrum is purely absolutely continuous. Under the assumption of flux rationality, the addition of a periodic potential does not alter the broadest features of $\sigma(H_{B,V})$; the spectrum remains purely absolutely continuous (this can be proved using the standard Thomas approach employed for $-\Delta + V$). Furthermore the spectrum contains at most finitely many gaps (the “Bethe-Sommerfeld conjecture”):

Theorem 1. *Suppose (FR) holds and V satisfies the regularity condition*

$$\sum_{m \in \mathbb{Z}^3} |m|^\delta |\widehat{V}_m| < +\infty$$

for some $\delta > 0$ (where \widehat{V}_m denote the Fourier coefficients of V). Then there exists $\Gamma \in \mathbb{R}$ such that $[\Gamma, \infty) \subseteq \sigma(H_{B,V})$; in particular, $\sigma(H_{B,V})$ contains only finitely many gaps. Furthermore, Γ depends continuously on the lattice Λ and on $\beta = |B|$.

See [1]; previously the result was obtained for any sufficiently small bounded V in [3].

There do not appear to be any general results on $\sigma(H_{B,V})$ in the case of non-rational flux.

3. DIMENSION $d = 2$

It is well known that $\sigma(H_{B,0})$ consists of the discrete eigenvalues $\beta(2n - 1)$, $n \in \mathbb{N}$, each of which has infinite multiplicity (the Landau levels). The presence of a non-zero potential V smears the Landau level $\beta(2n - 1)$ into a region of spectrum contained within the interval

$$I_{\beta,V}^n = \beta(2n - 1) + \widehat{V}_0 + C_{\beta,V} n^{-1/4}[-1, 1].$$

In particular, $\sigma(H_{B,V})$ contains infinitely many gaps for any $\beta \neq 0$ and V .

The character of the spectrum $\sigma(H_{B,V}) \cap I_{\beta,V}^n$ depends critically on the rationality of the flux. Under condition (FR), the form of the direct integral (2) makes it clear that (for sufficiently large n) $\sigma(H_{B,V}) \cap I_{\beta,V}^n$ will consist of p (possibly overlapping and/or degenerate) bands. The existence of eigenvalues for $V \not\equiv 0$ has not been fully resolved, although it appears the spectral bands are non-degenerate at least for generic V ([5]).

The study of $\sigma(H_{B,V}) \cap I_{\beta,V}^n$ in the case of non-rational flux has been undertaken in various limiting regimes (strong and weak magnetic fields were considered in [4]); in particular, it has been found that, after suitable normalisation, the limiting spectrum $\sigma(H_{B,V}) \cap I_{\beta,V}^n$ can be described by a Harper type operator. In this line, the following new result has been obtained for the large energy limit $n \rightarrow \infty$.

Theorem 2. *For all $n \gg 1$ there exists a neighbourhood Ω_n of $C_{\beta,V}[-1, 1] \subset \mathbb{C}$ and a holomorphic family of oscillatory pseudo-differential operators on $L^2(\mathbb{R})$, $Q_n(\mu)$, $\mu \in \Omega_n$, such that $\lambda \in \sigma(H_{B,V}) \cap I_{\beta,V}^n$ iff $0 \in \sigma(Q_n(\mu))$, where $\lambda = \beta(2n - 1) + \widehat{V}_0 + n^{-1/4}\mu$. Furthermore, as $n \rightarrow \infty$*

$$Q_n(\mu) = \text{Op}^w(W_n(x, \xi/\beta)) - \mu + O(n^{-1/4} \ln n),$$

where W_n is the periodic function given by

$$W_n(x, \xi) = \frac{(2\beta)^{1/4}}{\sqrt{\pi}} \sum_{m \in \mathbb{Z}^2 \setminus \{0\}} e^{i(m_1 x + m_2 \xi)} \frac{\widehat{V}_m}{\sqrt{|m|}} \cos\left(\frac{\sqrt{2}|m|}{\sqrt{\beta}} \sqrt{n} - \frac{\pi}{4}\right).$$

In particular, in the limit $n \rightarrow \infty$ the normalised spectrum $\sigma(H_{B,V}) \cap I_{\beta,V}^n$ is given as the spectrum of the operator $\text{Op}^w(W_n(x, \xi/\beta))$. This operator is in general of Harper type; in particular, for the potential $V(x, y) = \cos(x) + \cos(y)$ we get

$$\text{Op}^w(W_n(x, \xi/\beta)) = \frac{(2\beta)^{1/4}}{\sqrt{\pi}} \cos\left(\sqrt{2n/\beta} - \frac{\pi}{4}\right) (\cos(x) + \cos(D/\beta)),$$

which is a scaled version of the standard Harper operator (with parameter $1/\beta$).

In [4] it is shown that for certain V and large irrational magnetic fluxes β , the spectrum $\sigma(H_{B,V}) \cap I_{\beta,V}^n$ is a Cantor set (as is the case for the limiting Harper type operator). It is anticipated that similar results should be attainable for the large energy limit $n \rightarrow \infty$ (and probably for the weak electric field limit $V \rightarrow 0$).

The methods used to obtain Theorem 2 lead to an asymptotic formula for the eigenvalues of a harmonic oscillator perturbed by a (quasi-)periodic potential; these asymptotics are unusual in the sense that the leading order term contains an oscillatory factor, knowledge of which leads to the recovery of “half” the Fourier coefficients of V (see [2]).

REFERENCES

- [1] D. M. Elton, *The Bethe-Sommerfeld Conjecture for the Three-Dimensional Periodic Landau Operator*, to appear Rev. Math. Phys.
- [2] D. M. Elton, *Asymptotics for the Eigenvalues of the Harmonic Oscillator with a Quasi-Periodic Perturbation*, math.SP/0312110
- [3] V. A. Geĭler, V. A. Margulis and I. I. Chuchaev, *Spectrum Structure for the Three-Dimensional Periodic Landau Operator*, Algebra i Analiz **8** (1996), 104–124; English transl., St. Petersburg Math. J. **8** (1997), 447–461.
- [4] B. Helffer and J. Sjöstrand, *Equation de Schrödinger avec champ magnétique et équation de Harper*, 118–197, Lecture Notes in Phys., 345, Springer, Berlin, 1989.
- [5] F. Klopp, private communication.

Spectral Shift Function for Magnetic Schrödinger Operators

GEORGI RAIKOV

Let $H_0 := (i\nabla + A)^2 - b$ be the 3D magnetic Schrödinger operator essentially self-adjoint on $C_0^\infty(\mathbb{R}^3)$. Here $A = (-\frac{bx_2}{2}, \frac{bx_1}{2}, 0)$ is a magnetic potential which generates the constant magnetic field $B = \text{curl } A = (0, 0, b)$, $b > 0$. It is well-known that $\sigma(H_0) = \sigma_{\text{ac}}(H_0) = [0, \infty)$ (see e.g. [1]), where $\sigma(H_0)$ stands for the spectrum of H_0 , and $\sigma_{\text{ac}}(H_0)$ for its absolutely continuous spectrum. Moreover, the so-called Landau levels $2bq$, $q \in \mathbb{Z}_+$, play the role of thresholds in $\sigma(H_0)$.

Further, assume that the function V satisfies

$$(1) \quad V \not\equiv 0, \quad V \in C(\mathbb{R}^3), \quad 0 \leq V(\mathbf{x}) \leq c_0(1 + |\mathbf{x}|)^{-m}, \quad m > 3, \quad \mathbf{x} \in \mathbb{R}^3.$$

On the domain of H_0 define the operator $H_\pm := H_0 \pm V$ so that the electric potential $\pm V$ has a fixed sign. For every $E < \inf \sigma(H_\pm)$ we have $(H_\pm - E)^{-1} - (H_0 - E)^{-1} \in S_1$ where S_1 denotes the trace class. Hence, there exists a unique function $\xi = \xi(\cdot; H_\pm, H_0) \in L^1(\mathbb{R}; (1 + E^2)^{-1} dE)$ vanishing identically on $(-\infty, \inf \sigma(H_\pm))$, such that the *Lifshits-Krein trace formula*

$$\text{Tr}(f(H_\pm) - f(H_0)) = \int_{\mathbb{R}} \xi(E; H_\pm, H_0) f'(E) dE$$

holds for each $f \in C_0^\infty(\mathbb{R})$ (see [7, Chapter 8]). The function $\xi(\cdot; H_\pm, H_0)$ called the *spectral shift function* (SSF) for the operator pair (H_\pm, H_0) , is well defined on $\mathbb{R} \setminus 2b\mathbb{Z}_+$, bounded on every compact subset of $\mathbb{R} \setminus 2b\mathbb{Z}_+$, and continuous on

$\mathbb{R} \setminus \{2b\mathbb{Z}_+ \cup \sigma_{\text{pp}}(H_{\pm})\}$ where $\sigma_{\text{pp}}(H_{\pm})$ is the set of the eigenvalues of H_{\pm} (see [2]). In this talk based on the results of [3], we will discuss the asymptotic behaviour as $\lambda \rightarrow 0$ of $\xi(2bq + \lambda; H_{\pm}, H_0)$, the parameters $b > 0$ and $q \in \mathbb{Z}_+$ being fixed.

Let $h_0 := \left(i\frac{\partial}{\partial x_1} - \frac{bx_2}{2}\right)^2 + \left(i\frac{\partial}{\partial x_2} + \frac{bx_1}{2}\right)^2 - b$ be the Landau Hamiltonian essentially self-adjoint on $C_0^\infty(\mathbb{R}^2)$. It is well-known that $\sigma(h_0) = \cup_{q=0}^\infty \{2bq\}$, and each eigenvalue $2bq$, $q \in \mathbb{Z}_+$, has infinite multiplicity (see e.g. [1]). For $q \in \mathbb{Z}_+$ denote by $p_q = p_q(b)$ the orthogonal projection onto the eigenspace $\text{Ker}(h_0 - 2bq)$.

Assume that (1) holds. For $X_\perp := (x_1, x_2) \in \mathbb{R}^2$ set $W(X_\perp) := \int_{\mathbb{R}} V(X_\perp, x_3) dx_3$. Then the Toeplitz-type operator $p_q W p_q : L^2(\mathbb{R}^2) \rightarrow L^2(\mathbb{R}^2)$ satisfies $0 \leq p_q W p_q \in S_1$ and $\text{rank } p_q W p_q = \infty$ for each $q \in \mathbb{Z}_+$.

If $T = T^*$ is a compact operator, we denote by $n_+(s; T)$ the number of the eigenvalues of T lying on the interval (s, ∞) , $s > 0$, and counted with the multiplicities.

Theorem 1. [3, Theorem 3.1] *Assume that V satisfies (1). Fix $b > 0$ and $q \in \mathbb{Z}_+$. Then for each $\varepsilon \in (0, 1)$ we have*

$$(2) \quad \begin{aligned} \xi(2bq - \lambda; H_+, H_0) &= O(1), \quad \lambda \downarrow 0, \\ -n_+((1 - \varepsilon)2\sqrt{\lambda}; p_q W p_q) + O(1) &\leq \\ \xi(2bq - \lambda; H_-, H_0) &\leq \end{aligned}$$

$$(3) \quad -n_+((1 + \varepsilon)2\sqrt{\lambda}; p_q W p_q) + O(1), \quad \lambda \downarrow 0.$$

Estimate (2) shows that $\xi(2bq - \lambda; H_+, H_0)$ remains bounded while, since $\text{rank } p_q W p_q = \infty$, estimate (3) implies that $\xi(2bq - \lambda; H_-, H_0) \rightarrow -\infty$ as $\lambda \downarrow 0$. Suppose that V satisfies (1). For $\lambda \geq 0$ define the matrix-valued function

$$\mathbb{W}_\lambda = \mathbb{W}_\lambda(X_\perp) := \begin{pmatrix} w_{11} & w_{12} \\ w_{21} & w_{22} \end{pmatrix}, \quad X_\perp \in \mathbb{R}^2,$$

where

$$w_{11} := \int_{\mathbb{R}} V(X_\perp, x_3) \cos^2(\sqrt{\lambda}x_3) dx_3, \quad w_{22} := \int_{\mathbb{R}} V(X_\perp, x_3) \sin^2(\sqrt{\lambda}x_3) dx_3,$$

$$w_{12} = w_{21} := \int_{\mathbb{R}} V(X_\perp, x_3) \cos(\sqrt{\lambda}x_3) \sin(\sqrt{\lambda}x_3) dx_3.$$

Then the operator $p_q \mathbb{W}_\lambda p_q : L^2(\mathbb{R}^2)^2 \rightarrow L^2(\mathbb{R}^2)^2$ satisfies $0 \leq p_q \mathbb{W}_\lambda p_q \in S_1$ and $\text{rank } p_q \mathbb{W}_\lambda p_q = \infty$ for each $q \in \mathbb{Z}_+$ and $\lambda \geq 0$.

Theorem 2. [3, Theorem 3.2] *Assume that (1) holds. Fix $b > 0$ and $q \in \mathbb{Z}_+$. Then for each $\varepsilon \in (0, 1)$ we have*

$$\begin{aligned} \pm \frac{1}{\pi} \text{Tr} \arctan \left(((1 \pm \varepsilon)2\sqrt{\lambda})^{-1} p_q \mathbb{W}_\lambda p_q \right) + O(1) &\leq \\ \xi(2bq + \lambda; H_{\pm}, H_0) &\leq \\ \pm \frac{1}{\pi} \text{Tr} \arctan \left(((1 \mp \varepsilon)2\sqrt{\lambda})^{-1} p_q \mathbb{W}_\lambda p_q \right) + O(1), &\quad \lambda \downarrow 0. \end{aligned}$$

Since $\text{rank } p_q \mathbb{W}_\lambda p_q = \infty$, Theorem 2 implies that $\xi(2bq + \lambda; H_\pm, H_0) \rightarrow \pm\infty$ as $\lambda \downarrow 0$.

The main tool used in the proofs of Theorems 1 and 2 is the representation of the SSF due to A. Pushnitski (see [4]).

Combining Theorems 1 and 2 with some results on the eigenvalue asymptotics for compact Toeplitz-type operators obtained in [5] and [6], we can deduce more explicit asymptotic formulae describing the behaviour as $\lambda \rightarrow 0$ of $\xi(2bq + \lambda; H_\pm, H_0)$ under generic assumptions about the decay of the electric potential at infinity. Roughly speaking, these assumptions concern the cases where W admits a power-like decay at infinity, W decays exponentially, or the support of W is compact.

Corollary 3. [3, Corollaries 3.1, 3.2] *Let (1) hold. Fix $b > 0$ and $q \in \mathbb{Z}_+$.*

i) Assume that $W \in C^1(\mathbb{R}^2)$, and

$$W(X_\perp) = w_0(X_\perp/|X_\perp|)|X_\perp|^{-\alpha}(1 + o(1)), \quad |X_\perp| \rightarrow \infty,$$

$$|\nabla W(X_\perp)| \leq c_1(1 + |X_\perp|)^{-\alpha-1}, \quad X_\perp \in \mathbb{R}^2,$$

with $\alpha > 2$, $0 \leq w_0 \in C(\mathbb{S}^1)$, and $w_0 \not\equiv 0$. Then we have

$$\xi(2bq - \lambda; H_-, H_0) = -\psi_\alpha(2\sqrt{\lambda})(1 + o(1)), \quad \lambda \downarrow 0,$$

$$\xi(2bq + \lambda; H_\pm, H_0) = \pm \frac{1}{2 \cos(\pi/\alpha)} \psi_\alpha(2\sqrt{\lambda})(1 + o(1)), \quad \lambda \downarrow 0,$$

where $\psi_\alpha(s) := s^{-2/\alpha} \frac{b}{4\pi} \int_{\mathbb{S}^1} w_0(t)^{2/\alpha} dt$, $s > 0$.

ii) Assume that $W \in L^\infty(\mathbb{R}^2)$, and

$$\ln W(X_\perp) = -\mu|X_\perp|^{2\beta}(1 + o(1)), \quad |X_\perp| \rightarrow \infty,$$

with some $\mu > 0$, and $\beta > 0$. Suppose in addition that V satisfies the estimate

$$(4) \quad V(X_\perp, x_3) \leq c_2(1 + |X_\perp|)^{-m_\perp}(1 + |x_3|)^{-m_3}, \quad X_\perp \in \mathbb{R}^2, \quad x_3 \in \mathbb{R},$$

with $m_\perp > 2$, $m_3 > 2$. Then we have

$$\xi(2bq - \lambda; H_-, H_0) = -\varphi_\beta(2\sqrt{\lambda})(1 + o(1)), \quad \lambda \downarrow 0,$$

$$\xi(2bq + \lambda; H_\pm, H_0) = \pm \frac{1}{2} \varphi_\beta(2\sqrt{\lambda})(1 + o(1)), \quad \lambda \downarrow 0,$$

where

$$\varphi_\beta(s) := \begin{cases} \frac{b}{2\mu^{1/\beta}} |\ln s|^{1/\beta} & \text{if } 0 < \beta < 1, \\ \frac{1}{\ln(1+2\mu/b)} |\ln s| & \text{if } \beta = 1, \\ \frac{\beta}{\beta-1} (\ln |\ln s|)^{-1} |\ln s| & \text{if } 1 < \beta < \infty, \end{cases} \quad s \in (0, e^{-1}).$$

iii) Finally, assume that $W \in L^\infty(\mathbb{R}^2)$, $\text{supp } W$ is compact, and there exists a constant $c > 0$ such that $W \geq c$ on an open non-empty subset of \mathbb{R}^2 . Suppose in addition that V satisfies (4) with $m_\perp > 2$, $m_3 > 2$. Then we have

$$\xi(2bq - \lambda; H_-, H_0) = -\varphi_\infty(2\sqrt{\lambda})(1 + o(1)), \quad \lambda \downarrow 0,$$

$$\xi(2bq + \lambda; H_\pm, H_0) = \pm \frac{1}{2} \varphi_\infty(2\sqrt{\lambda})(1 + o(1)), \quad \lambda \downarrow 0,$$

where

$$\varphi_\infty(s) := (\ln |\ln s|)^{-1} |\ln s|, \quad s \in (0, e^{-1}).$$

REFERENCES

- [1] J. AVRON, I. HERBST, B. SIMON, *Schrödinger operators with magnetic fields. I. General interactions*, Duke Math. J. **45** (1978), 847-883.
- [2] V. BRUNEAU, A. PUSHNITSKI, G. D. RAIKOV, *Spectral shift function in strong magnetic fields*, Algebra i Analiz **16** (2004), 207 - 238.
- [3] C. FERNÁNDEZ, G. D. RAIKOV, *On the singularities of the magnetic spectral shift function at the Landau levels*, **5** (2004), 381 - 403.
- [4] A. PUSHNITSKIĬ, *A representation for the spectral shift function in the case of perturbations of fixed sign*, Algebra i Analiz **9** (1997), 197-213 [in Russian]; English translation in St. Petersburg Math. J. **9** (1998), 1181-1194.
- [5] G. D. RAIKOV, *Eigenvalue asymptotics for the Schrödinger operator with homogeneous magnetic potential and decreasing electric potential. I. Behaviour near the essential spectrum tips*, Commun. P.D.E. **15** (1990), 407-434.
- [6] G.D. RAIKOV, S. WARZEL, *Quasi-classical versus non-classical spectral asymptotics for magnetic Schrödinger operators with decreasing electric potentials*, Rev. Math. Phys. **14** (2002), 1051-1072.
- [7] D. R. YAFAEV, *Mathematical scattering theory. General theory*. Translations of Mathematical Monographs, **105** AMS, Providence, RI, 1992.

On the Laplacian in the halfspace with a periodic boundary condition

RUPERT L. FRANK

The characteristic feature of Schrödinger operators that are periodic with respect to some, but not all directions is the appearance of surface states, see [1] and the references in [5], [6]. On physical grounds one expects that these states are not bound but correspond to additional channels of scattering, i.e., that the spectrum of the corresponding operator is purely absolutely continuous. We are only aware of [2], [3], [4], [5] dealing with this problem.

Here we follow [4] and study spectral and scattering properties of the Laplacian

$$H^{(\sigma)}u = -\Delta u \quad \text{on } \mathbb{R}_+^{d+1} := \{(x, y) \in \mathbb{R}^d \times \mathbb{R} : y > 0\}$$

together with a boundary condition of the third type

$$\frac{\partial u}{\partial \nu} + \sigma u = 0 \quad \text{on } \mathbb{R}^d \times \{0\}$$

with a $(2\pi\mathbb{Z})^d$ -periodic function $\sigma : \mathbb{R}^d \rightarrow \mathbb{R}$. Under the condition

$$(1) \quad \sigma \in L_{q,loc}(\mathbb{R}), \quad q > 1, \quad \text{if } d = 1, \quad \sigma \in L_{d,loc}(\mathbb{R}^d) \text{ if } d \geq 2,$$

$H^{(\sigma)}$ can be defined as a self-adjoint operator in $L_2(\mathbb{R}_+^{d+1})$ by means of the lower semibounded and closed quadratic form

$$\int_{\mathbb{R}_+^{d+1}} |\nabla u(x, y)|^2 dx dy + \int_{\mathbb{R}^d} \sigma(x) |u(x, 0)|^2 dx, \quad u \in H^1(\mathbb{R}_+^{d+1}).$$

Note that $H^{(\sigma)}$ can be viewed as a Schrödinger-type operator with singular potential $\sigma(x)\delta(y)$ describing the interaction of a quantum-mechanical particle with the surface of a crystal.

We investigate the scattering with respect to the Neumann Laplacian $H^{(0)}$.

Theorem 1. *Assume that σ satisfies (1). Then the wave operators*

$$W_{\pm}^{(\sigma)} := s - \lim_{t \rightarrow \pm\infty} \exp(itH^{(\sigma)}) \exp(-itH^{(0)})$$

exist and satisfy $\mathcal{R}(W_+^{(\sigma)}) = \mathcal{R}(W_-^{(\sigma)})$.

If σ is non-negative, we obtain a rather complete result.

Theorem 2. *Assume that σ satisfies (1) and $\sigma(x) \geq 0$ for a.e. $x \in \mathbb{R}^d$. Then the wave operators $W_{\pm}^{(\sigma)}$ are unitary and satisfy $H^{(\sigma)} = W_{\pm}^{(\sigma)} H^{(0)} W_{\pm}^{(\sigma)*}$.*

However, in general the wave operators will *not* be complete due to the existence of *surface states*, i.e., states that are localized near the boundary for all time. These states correspond to bands in the spectrum of $H^{(\sigma)}$. A sufficient condition for $\sigma(H^{(\sigma)}(k)) \cap (-\infty, 0) \neq \emptyset$ is

$$\int_{(-\pi, \pi)^d} \sigma(x) dx \leq 0, \quad \sigma \not\equiv 0.$$

It is natural to ask whether the spectrum of $H^{(\sigma)}$ is still absolutely continuous in this situation.

Theorem 3. *Assume that σ satisfies (1) if $d \leq 4$ and $\sigma \in L_{2(d-2),loc}(\mathbb{R}^d)$ if $d \geq 5$. Then the operator $H^{(\sigma)}$ has purely absolutely continuous spectrum.*

Hence surface states correspond to *additional channels of scattering*.

Let us explain some of the mathematical ideas involved. By means of Floquet theory we represent $H^{(\sigma)}$ as a direct integral

$$\int_{[-\frac{1}{2}, \frac{1}{2}]^d} \oplus H^{(\sigma)}(k) dk$$

with operators $H^{(\sigma)}(k)$ acting in $L_2(\Pi)$ where $\Pi := (-\pi, \pi)^d \times \mathbb{R}_+$ is a halfcylinder. The investigation of the operator $H^{(\sigma)}$ reduces to the study of the fibers $H^{(\sigma)}(k)$. Note that the fundamental domain Π is unbounded, so the operators $H^{(\sigma)}(k)$ have continuous spectrum. This part can be studied by scattering theory. To prove the absolute continuity of the spectrum of $H^{(\sigma)}$ we cannot (directly) apply the Thomas approach, since eigenvalues of $H^{(\sigma)}(k)$ may be embedded in the

continuous spectrum. We "separate" them from the remaining spectrum by characterizing them, in the spirit of the Birman-Schwinger principle, as parameters λ for which a pseudo-differential operator $B^{(\sigma)}(\lambda, k)$ on the boundary $(-\pi, \pi)^d \times \{0\}$ has eigenvalue 0. The latter operator has discrete spectrum and can be handled by Thomas' method.

Acknowledgements. The author is deeply grateful to Prof. M. Sh. Birman for the setting of the problem, useful discussions and constant attention to the work. It is a great pleasure to thank R. G. Shterenberg for numerous consultations and Prof. A. Sobolev and Prof. T. Weidl for the invitation to the wonderful conference.

REFERENCES

- [1] E. B. Davies, B. Simon, *Scattering Theory for Systems with Different Spatial Asymptotics on the Left and Right*, Commun. Math. Phys. **63** (1978), 277-301.
- [2] N. Filonov, F. Klopp, *Absolute continuity of the spectrum of a Schrödinger operator with a potential which is periodic in some directions and decays in others*, Documenta Math. **9** (2004), 107-121; Erratum: ibd., 135-136.
- [3] N. Filonov, F. Klopp, *Absolutely continuous spectrum for the isotropic Maxwell operator with coefficients that are periodic in some directions and decay in others*, Comm. Math. Phys., to appear.
- [4] R. L. Frank, *On the Laplacian in the halfspace with a periodic boundary condition*, preprint, mp_arc 04-407.
- [5] R. L. Frank, R. G. Shterenberg, *On the scattering theory of the Laplacian with a periodic boundary condition. II. Additional channels of scattering*, Documenta Math. **9** (2004), 57-77.
- [6] P. Kuchment, *On some spectral problems of mathematical physics*, Partial differential equations and inverse problems, 241-276, Contemp. Math., 362, Amer. Math. Soc., Providence, RI, 2004.

A Multidimensional Trace Formula

A. LAPTEV

(joint work with S. Naboko and O. Safronov)

Let us consider the equation

$$(1) \quad Hu = -\Delta u + Vu = k^2 u,$$

where $V \in C_0^\infty(\mathbb{R}^3)$ and $\text{supp} V \subset \{x : c_1 < |x| < c_2\}$, $c_1, c_2 > 0$. By using the unitary transformation U from $L^2((0, \infty), dr; L^2(\mathbb{S}^2))$ to $L^2((0, \infty), r^2 dr; L^2(\mathbb{S}^2))$,

$$v(t, \theta) = Uu(r, \theta) = r^{-1}u,$$

we reduce the study of (1) the operator \tilde{H} in $L^2((0, \infty), dr; L^2(\mathbb{S}^2))$

$$(2) \quad \tilde{H}v = -\partial_{rr}^2 v + \frac{B}{r^2}v + Vv = k^2 v,$$

where B is the Laplace-Beltrami operator in $L^2(\mathbb{S}^2)$.

We now consider the equation

$$(3) \quad -f_{rr}''(r, \theta, k) + \frac{B}{r^2} f(r, \theta, k) + Vf(r, \theta, k) = k^2 f(r, \theta, k),$$

subject to the initial condition

$$(4) \quad f(r, \theta, k) = e^{-ikr}, \quad 0 < r < c_1.$$

It can be shown that there are "scattering" coefficients a and b such that

$$(5) \quad f(r, \theta, k) = a(\theta, k)e^{-ikr} \left(1 + O(|kr|^{-1})\right) + b(\theta, k)e^{ikr} \left(1 + O(|kr|^{-1})\right), \quad r \rightarrow \infty.$$

If $f(r, \theta, k)$, $k \in \mathbb{C}$, $r \geq 1$, is a solution of the differential equation (3) then it satisfies the integral equation

$$(6) \quad f(r, \theta, k) = e^{-ikr} - \frac{1}{2ik} \int_0^r \left(e^{-ik(r-t)} - e^{ik(r-t)}\right) \left(V(t, \theta) + \frac{B}{r^2}\right) f(t, \theta, k) dt.$$

Substituting $\psi(r, \theta, k) = e^{ikr} f(r, \theta, k)$ we obtain

$$(7) \quad \psi(r, \theta, k) = 1 - \mathcal{K} \psi(r, \theta, k) = \int_0^r K(r, t, k) \psi(t, \theta, k) dt,$$

where by \mathcal{K} we denote the integral operator whose operator valued symbol is equal to

$$(8) \quad K(r, t, k) = \frac{(1 - e^{2ik(r-t)})}{2ik} \left(V(t, \cdot) + \frac{B}{r^2}\right).$$

Solving the Volterra equation (6) we obtain the series

$$\psi(r, \theta, k) = 1 + \sum_{j=1}^{\infty} \int_{r \geq t_1 \geq \dots \geq t_m \geq 0} \prod_{q=1}^j K(t_{l-1}, t_l, k) dx_1 \cdots dx_j \cdot 1.$$

This series is convergent pointwise and, in particular, $\psi(r, \theta, k) \equiv 1$ if $0 \leq r \leq c_1$. The function $\psi(r, \theta, k)$ is smooth and also analytic with respect to $k \in \mathbb{C} \setminus \{0\}$. Indeed, since the kernel $K(r, t, k)$ is analytic in k , we obtain

$$\frac{\partial}{\partial \bar{k}} \psi(r, \theta, k) = - \int_0^r K(r, t, k) \frac{\partial}{\partial \bar{k}} \psi(t, \theta, k) dt.$$

Therefore $\partial \psi(r, \theta, k) / \partial \bar{k}$ satisfies a homogeneous Volterra integral equation and hence it is identically equal to zero.

The Volterra equation (6) can be rewritten as

$$(9) \quad f(r, \theta, k) = e^{-ikr} \left[1 - \frac{1}{2ik} \int_0^r V(t, \theta, k) dt - \frac{1}{2ik} \int_0^r \left(V(t, \theta, k) + \frac{B}{r^2}\right) (\psi(t, \theta) - 1) dt \right] + \frac{e^{ikr}}{2ik} \left[\int_0^r e^{-2ikt} V(t, \theta, k) dt + \int_0^r e^{-2ikt} \left(V(t, \theta, k) + \frac{B}{r^2}\right) (\psi(t, \theta) - 1) dt \right].$$

Comparing (9) with (5) we see that

$$(10) \quad a(\theta, k) = 1 - \frac{1}{2ik} \int_0^r V(t, \theta, k) dt$$

$$\begin{aligned}
& -\frac{1}{2ik} \int_0^r \left(V(t, \theta, k) + \frac{B}{r^2} \right) (\psi(t, \theta) - 1) dt \\
(11) \quad & b(\theta, k) = \frac{1}{2ik} \int_0^r e^{-2ikt} V(t, \theta, k) dt \\
& + \frac{1}{2ik} \int_0^r e^{-2ikt} \left(V(t, \theta, k) + \frac{B}{r^2} \right) (\psi(t, \theta) - 1) dt.
\end{aligned}$$

Note that for a fixed k_0 , $\text{Im } k_0 > 0$, if we assume that the function f is from the class $L^2((0, \infty) \times \mathbb{S}^2)$, then $a(\theta, k_0)$ is equal to zero identically in θ . This implies that $a(\theta, k_0) \equiv 0$ if k_0 is an eigenvalue of the operator (1).

Let \varkappa_j , $j = 1, \dots, J$, be zeros of the function $\int_{\mathbb{S}^2} a(\theta, k) d\theta$ in the upper half plane. We obtain a version of Buslaev-Faddeev-Zakharov trace formula, see [1] and [2].

Theorem 1. *Let V be a $C_0^\infty(\mathbb{R}^3)$ and $\text{supp } V \subset \{x : c_1 < |x| < c_2\}$, $c_1, c_2 > 0$. Then the following trace formula holds true*

$$\begin{aligned}
& \sum_j \varkappa_j^3 + \frac{3}{2\pi} \int_{-\infty}^{\infty} k^2 \log \left| \int_{\mathbb{S}^2} a(\theta, k) d\theta \right| dk \\
& = \frac{3}{16} \int_{\mathbb{S}^2} \int_0^\infty \left\{ \left| \int_0^r \nabla_\theta V(t, \theta) dt \right|^2 r^{-2} + V^2(r, \theta) \right\} dr d\theta.
\end{aligned}$$

When proving the theorem we use an approach developed in [3], where the authors have considered trace formulae with operator valued potentials and their applications. Similar ideas have been also used in [4] when proving absolute continuity of the spectrum of Schrödinger operators with oscillating potentials.

Acknowledgements. A.L and O.S. thank a partial support by the SPECT ESF European programme. S.N. was also partly supported by the KBN grant 5, PO3A/026/21. g1925l.

REFERENCES

- [1] V.S Buslaev and L.D. Faddeev, *Formulas for traces for a singular Sturm-Liouville differential operator*. [English translation], Dokl. AN SSSR, **132** (1960), 451-454.
- [2] L.D Faddeev and V.E. Zakharov, *Korteweg-de Vries equation: A completely integrable hamiltonian system*. Func. Anal. Appl., **5** (1971), 18-27.
- [3] A. Laptev and T. Weidl, *Sharp Lieb-Thirring inequalities in high dimensions*, Acta Mathematica **184** (2000), 87-111.
- [4] A. Laptev, S. Naboko and O. Safronov, *A Szegö condition for a multidimensional Schrödinger operator*, J.Func.Anal., **219** (2005), 285-305.

Participants

Dr. Jean-Marie Barbaroux

Departement de Mathematiques
Universite de Toulon et du Var
B.P. 132
F-83957 La Garde Cedex

Prof. Dr. Rafael Benguria

Facultad de Fisica
Pontificia Univ. Catolica
Casilla 306
Santiago 22 – Chile

Prof. Dr. Michiel van den Berg

School of Mathematics
University of Bristol
University Walk
GB-Bristol BS8 1TW

Prof. Dr. Mikhail S. Birman

Dept. of Mathematical Physics
Institute of Physics
St. Petersburg State University
Petrodvoretz, Ulyanov St. 1
198904 St. Petersburg – Russia

Dr. David Damanik

Department of Mathematics
253-37
California Institute of Technology
Pasadena, CA 91125 – USA

Prof. Dr. Tomas Ekholm

Dept. of Mathematics
Royal Institute of Technology
Lindstedtsvägen 25
S-100 44 Stockholm

Dr. Daniel M. Elton

Dept. of Mathematics & Statistics
University of Lancaster
Fylde College
Bailrigg
GB-Lancaster, LA1 4YF

Prof. Dr. Laszlo Erdős

Mathematisches Institut
Universität München
Theresienstr. 39
D-80333 München

Prof. Dr. William Desmond Evans

School of Mathematics
Cardiff University
23, Senghennydd Road
GB-Cardiff CF24 4AG

Prof. Dr. Pavel Exner

Department of Theoretical Physics
Nuclear Physics Inst.
Academy of Sciences
25068 Rez (near Prague) – Czech Republ.

Clemens Förster

Inst. f. Analysis, Dynamik, Modellierung
Universität Stuttgart
Pfaffenwaldring 57
D-70569 Stuttgart

Dipl.-Math. Rupert L. Frank

Department of Mathematics
Royal Institute of Technology
S-10044 Stockholm

Prof. Dr. Friedrich Götze

Fakultät für Mathematik
Universität Bielefeld
Postfach 100131
D-33501 Bielefeld

Dr. Gian Michele Graf

Institut für Theoretische Physik
ETH Zürich
Hönggerberg
CH-8093 Zürich

Prof. Dr. Bernard Helffer

Department of Mathematics
Univ. Paris-Sud
Bat. 425
F-91405 Orsay Cedex

Prof. Dr. Rainer Hempel

Institut Computational Mathematics
Technische Universität Braunschweig
Pockelstraße 14
D-38106 Braunschweig

Prof. Dr. Ira W. Herbst

Dept. of Mathematics
University of Virginia
Kerchof Hall
P.O.Box 400137
Charlottesville, VA 22904-4137 – USA

Prof. Dr. Peter David Hislop

Dept. of Mathematics
University of Kentucky
Lexington, KY 40506-0027 – USA

Prof. Dr. M. Hoffmann-Ostenhof

Fakultät für Mathematik
Universität Wien
Nordbergstr. 15
A-1090 Wien

Prof. Dr. Th. Hoffmann-Ostenhof

Institut für Theoretische Chemie
Universität Wien
Währingstr. 17
A-1090 Wien

Prof. Dr. Dirk Hundertmark

Dept. of Mathematics, University of
Illinois at Urbana-Champaign
273 Altgeld Hall MC-382
1409 West Green Street
Urbana, IL 61801-2975 – USA

Prof. Dr. Arne Jensen

Dept. of Mathematical Sciences
University of Aalborg
Fredrik Bajers Vej 7G
DK-9220 Aalborg East

Prof. Dr. Werner Kirsch

Mathematisches Institut
Ruhr Universität Bochum
Universitätsstraße 150
D-44780 Bochum

Prof. Dr. Alexander Kiselev

Department of Mathematics
University of Wisconsin-Madison
480 Lincoln Drive
Madison, WI 53706-1388 – USA

Prof. Dr. Frederic Klopp

Departement de Mathematiques
Institut Galilee
Universite Paris XIII
99 Av. J.-B. Clement
F-93430 Villetaneuse

Dr. Vadim Kostrykin

Fraunhofer-Institut
für Lasertechnik
Steinbachstr. 15
D-52074 Aachen

Dr. Hynek Kovarik

Institut für Analysis, Dynamik
und Modellierung
Universität Stuttgart
Pfaffenwaldring 57
D-70569 Stuttgart

Prof. Dr. Peter Kuchment

Department of Mathematics
Texas A & M University
College Station, TX 77843-3368 – USA

Prof. Dr. Ari Laptev

Dept. of Mathematics
Royal Institute of Technology
Lindstedtsvägen 25
S-100 44 Stockholm

Dipl. Math. Helmut Linde

Facultad de Fisica
Pontificia Univ. Catolica
Casilla 306
Santiago 22 – Chile

Prof. Dr. Shu Nakamura

Graduate School of
Mathematical Sciences
University of Tokyo
3-8-1 Komaba, Meguro-ku
Tokyo 153-8914 – JAPAN

Dr. Yuri Netrusov

Department of Mathematics
University of Bristol
University Walk
GB-Bristol, BS8 1TW

Dr. Leonid Parnovski

Department of Mathematics
University College London
Gower Street
GB-London, WC1E 6BT

Dr. Alexander Pushnitski

Department of Mathematical Sciences
Loughborough University
Loughborough
GB-Leicestershire LE11 3TU

Prof. Dr. Georgi D. Raikov

Departamento de Matematicas
Universidad de Chile
Las Palmaras 3425
Casilla 653
Santiago – Chile

Prof. Dr. Didier Robert

Laboratoire de Mathematiques
Jean Leray UMR 6629
Universite de Nantes, B.P. 92208
2 Rue de la Houssiniere
F-44322 Nantes Cedex 03

Dr. Norbert Röhrl

Institut für Analysis, Dynamik
und Modellierung
Universität Stuttgart
Pfaffenwaldring 57
D-70569 Stuttgart

Dr. Roman G. Shterenberg

Dept. of Mathematical Physics
Institute of Physics
St. Petersburg State University
Petrodvoretz, Ulyanov St. 1
198904 St. Petersburg – Russia

Prof. Dr. Heinz Karl H. Siedentop

Mathematisches Institut
Universität München
Theresienstr. 39
D-80333 München

Prof. Dr. Maxim Skriganov

St. Petersburg Branch of
Mathematical Institute of
Russian Academy of Science
Fontanka 27
191023 St. Petersburg – Russia

Prof. Dr. Alexander V. Sobolev

Department of Mathematics
University of Sussex
Mantell Bldg.
Falmer
GB-Brighton BN1 9RF

Prof. Dr. Michael Solomyak

Department of Mathematics
The Weizmann Institute of Science
P. O. Box 26
Rehovot 76 100 – Israel

Dr. Edgardo Stockmeyer

Mathematisches Institut
Universität München
Theresienstr. 39
D-80333 München

Prof. Dr. Tatyana Suslina

Dept. of Mathematical Physics
Institute of Physics
St. Petersburg State University
Petrodvoretz, Ulyanov St. 1
198904 St. Petersburg – Russia

Prof. Dr. Timo Weidl

Institut für Analysis, Dynamik
und Modellierung
Universität Stuttgart
Pfaffenwaldring 57
D-70569 Stuttgart

Dr. Semjon Wugalter

Mathematisches Institut
Universität München
Theresienstr. 39
D-80333 München

Prof. Dr. Dimitrij Rael Yafaev

U. F. R. Mathematiques
I. R. M. A. R.
Universite de Rennes I
Campus de Beaulieu
F-35042 Rennes Cedex

Prof. Dr. Kenji Yajima

Department of Mathematics
Faculty of Science
Gakushuin University
Toshimaku-ku, 1-5-1 Mejiro
Tokyo 171-8588 – JAPAN

