

Some Rigorous Results for the Anderson Model

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We discuss two results for the Anderson model of random quantum Hamiltonians: (1) smoothness of the density of states in the one-dimensional model, even in many cases where the potential distribution is not smooth; and (2) a criterion for localization which, among other consequences, implies that certain estimates of Fröhlich and Spencer yield a dense point spectrum for the multidimensional model at large randomness or large energies.

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In this Letter, we want to announce some rigorous results, and discuss their consequences on the regularity of the density of states in the Anderson model obtained by two of us¹, and on the localization in this model obtained by two of us.² Full details will appear elsewhere.^{1,2}

The Anderson model of random impurities³ is the random Hamiltonian $H_\omega = H_0 + V_\omega$ on $l^2(Z^\nu)$, where

$$(H_0 u)(n) = \sum_{|j|=1} u(n+j)$$

and V_ω is the diagonal operator $V_\omega(n)$, with $V_\omega(n)$ independent identically distributed random variables with distribution $d\kappa(v)$.

Theorem 1 (Ref. 1).—In the one-dimensional case ($\nu = 1$), suppose that $d\kappa$ has the form $d\kappa(v) = F(v)dv$, where F has compact support and $\hat{F}(k) \equiv \int e^{-ikv} F(v) dv$ obeys⁴

$$|\hat{F}(k)| \leq C(1 + |k|)^{-\alpha}, \quad \alpha > \frac{1}{2}.$$

Then, the integrated density of states, $k(E)$, is an infinitely differentiable function.

This result, whose proof we discuss at the end of this Letter, says that $k(E)$ can be much smoother than the distribution of $V(n)$. Previous results either proved some form of continuity weaker than differentiability⁵ or proved that k is only at least as smooth as $\int_{-\infty}^E d\kappa(v)$.⁶ The result applies to the case where $V(n)$ is uniformly distributed in some interval $[a, b]$. Infinite smoothness of $k(E)$ is consistent with, and suggested by, the phenomenon of Lifshitz tails.

Halperin⁷ has proven that when $d\kappa = \theta\delta(v-a) + (1-\theta)\delta(v-b)$ and either $|a-b|$ is large or $\theta(1-\theta)$ is small, then k is not differentiable; indeed, it is not Hölder continuous of any prescribed order.

This shows that some hypothesis on $d\kappa$ is needed.

The second result is a criterion for localization; $G_\omega(n, m; z)$ is the Green's function $(\delta_m, (H_\omega - z)^{-1} \delta_n)$:

Theorem 2 (Ref. 2).—Suppose that for almost all $E \in (a, b)$ and almost all ω we have that

$$\sup_{0 < \epsilon < 1} \left| \sum_n |G_\omega(0, n; E + i\epsilon)|^2 \right| < \infty. \quad (1)$$

If $d\kappa$ has an absolutely continuous component⁸ and $\nu = 1$, or $d\kappa$ is absolutely continuous and ν is arbitrary, then for almost all ω , H_ω has only a point spectrum⁹ in (a, b) . If the essential support of the absolutely continuous component of $d\kappa$ is $(-\infty, \infty)$,¹⁰ then (1) is not only sufficient for a pure point spectrum, it is also necessary.

The quantity on the lefthand side of (1) increases as ϵ decreases, and so we need only treat sufficiently small ϵ . One estimate which clearly implies (1) is

$$|G_\omega(0, n; E + i\epsilon)| \leq C_{\omega, E} \exp[-C(E)|n|] \quad (2)$$

for almost every (a.e.) ω and all sufficiently small ϵ . In this case, one can prove² that the eigenfunctions decay with a localization length¹¹ no larger than $C(E)^{-1}$ so long as $d\kappa$ is purely absolutely continuous.

Before discussing the proof of theorem 2, we note that there are two cases where one knows how to prove (1) [in fact, to prove (2)]: in the general one-dimensional case, and in the higher-dimensional case at strong coupling. In the one-dimensional case, Ishii and Deift-Simon^{12,13} proved (2). This provides a new proof of localization in this case. The point is not so much that our hypothesis on $d\kappa$ is weaker than that in existing proofs,¹⁴ but that the proof via theorem 2 is

mathematically and conceptually quite simple, and more significantly, explains why (2) can hold in certain almost-periodic models¹⁵ which only have singular continuous spectra.¹⁶ Ishii's bounds, together with general lower bounds on eigenfunctions,¹⁷ imply that the localization length is the inverse of the Lyapunov exponent.¹⁸

Fröhlich and Spencer¹⁹ have proven (2) in the multidimensional Anderson model under two circumstances: (i) $d\kappa$ Gaussian and $|E|$ very large, and (ii) $d\kappa = g(E) dE$ with $\sup_E |g(E)|$ sufficiently small (large coupling or large randomness). While it was known that these estimates imply the absence of extended states,²⁰ it was not known until now that the estimates of Ref. 19 imply a point spectrum. Recently, Fröhlich *et al.*²¹ and Goldsheid²² have announced results on localization in the multidimensional situations. The Fröhlich-Spencer estimates¹⁹ and our remarks on the localization length imply that the localization length goes to zero in the infinite-randomness or large-energy limit.

Theorem 2 comes from an analysis of the spectrum of self-adjoint operators under a random rank-one perturbation. The basic deterministic theory of such perturbations was developed by Aronszajn²³ and Donoghue,²⁴ and our own interest was kindled by the recent work of Kotani²⁵ on the special case of random boundary conditions in half-line problems. Indeed, the proof of theorem 3 below is essentially a synthesis of ideas of Aronszajn and Kotani.

Let A be a self-adjoint operator, let P be the projection onto a unit vector, ϕ , and let $A_\lambda = A + \lambda P$. Let $d\mu_\lambda$ be the spectral measure²⁶ defined by

$$(\phi, e^{-iA_\lambda} \phi) = \int e^{-ix} d\mu_\lambda(x).$$

We need two functions related to these measures:

$$F_\lambda(z) = \int (x-z)^{-1} d\mu_\lambda(x), \tag{3}$$

$$B(x) = \left[\int (x-y)^{-2} d\mu_0(x) \right]^{-1}.$$

The Steiltjes transform, $F_\lambda(z)$, is analytic in the upper half-plane, and the general theory of such functions²⁷ implies that boundary values $F_\lambda(x+i0)$ exist (for λ fixed) for almost all x . Since $\text{Im}F_0(x+i\epsilon) \leq (\text{Im}\epsilon)B(x)^{-1}$, at most one of $\text{Im}F_0(x+i0)$ and $B(x)$ is nonzero at any point.

Theorem 3. $-d\mu_\lambda$ has a vanishing singular continuous part for almost all λ if and only if $B(x) + \text{Im}F_0(x+i0) > 0$ for almost all x .

Before discussing the proof of this theorem, we explain how it implies theorem 2. If $d\mu_0$ is the spectral measure for H_ω associated to δ_0 , then a simple calculation shows that the left-hand side of Eq. (1) is $B(E)^{-1}$, and so (1) says that for a.e. ω and a.e.

$E \in (a,b)$, $B(E) > 0$. As noted above, this implies that $\text{Im}F_0(E+i0) = 0$, and thus by Eq. (4) below, $\text{Im}F_\lambda(E+i0) = 0$. The general theory of boundary values of Steiltjes transforms²⁷ implies that $d\mu_\lambda^{\text{sc}} = 0$ on (a,b) . Thus, theorem 3 says that for a.e. λ , $H_\omega + \lambda P_0$ has only a point spectrum for a.e. ω and a.e. λ . The λP_0 just shifts the value of $V(0)$. Since $V(0)$ is independent of the other $V(n)$'s and $d\kappa$ has an absolutely continuous component, we have a point spectrum in the original H_ω with nonzero probability, and so with probability 1 by general results (see the first reference in Ref. 15).

Here is a sketch of the proof of theorem 3: (i) By taking expectations in $(A_\lambda - z)^{-1} = (A_0 - z)^{-1} - \lambda(A_0 - z)^{-1}P(A_\lambda - z)^{-1}$, we obtain the basic equation of Aronszajn²³:

$$F_\lambda(z) = F_0(z) / [1 + \lambda F_0(z)]. \tag{4}$$

(ii) Since $\mu_\lambda(\{E_0\}) = \lim_{\epsilon \downarrow 0} i\epsilon F_\lambda(E_0 + i\epsilon)$, one deduces from (4) that $\mu_\lambda(\{E_0\}) > 0$ if and only if $F_0(E_0 + i0) = -\lambda^{-1}$ and $B(E_0) > 0$; in fact, $\mu_\lambda(\{E_0\}) = \lambda^{-2}B(E_0)$ if $F_0(E_0 + i0) = -\lambda^{-1}$. (iii) By using Eq. (4), one can study the measure $d\eta$ defined by

$$\int g(E) d\eta(E) = \int \left[\int g(E) d\mu_\lambda(E) \right] (1 + \lambda^2)^{-1} d\lambda.$$

From (4), one finds that

$$F^{(\eta)}(z) \equiv \int (x-z)^{-1} d\eta(x) = \pi / [F_0(z)^{-1} - i]. \tag{5}$$

From this, one can deduce that $d\eta(x) = H(x) dx$, where H is almost everywhere nonzero. For example, since $\text{Im}F_0(z) > 0$, $\text{Im}F^\eta(z) \leq \pi$, which implies that $H(x) = \pi^{-1} \text{Im}F^\eta(x+i0) \leq 1$. (iv) Let $C = \{x | F_0(x+i0) = -\lambda^{-1}; B(x) = 0\}$. The theorem of de Vallée Poussain²⁷ says that the singular continuous part of $d\mu_\lambda$, call it $d\mu_\lambda^{\text{sc}}$, is supported on the set where $F_\lambda(x+i0) = \infty$, which, by (4), is the set of where $F_0(x+i0) = -\lambda^{-1}$. By (ii), the subset of this set where $B(x) > 0$ consists of point masses of $d\mu_\lambda$ so that it is countable. Thus $\mu_\lambda(C) = \mu_\lambda^{\text{sc}}(C)$. (v) By (iii), $|C| = \int_C dx = 0$ if and only if $\int \mu_\lambda(C) (1 + \lambda^2)^{-1} d\lambda = 0$ which, by (iv), is true if and only if $\mu_\lambda^{\text{sc}}(C) = 0$ for almost all λ . This completes the proof of theorem 3.

In some ways, the key step is (iii) (related to Kotani's work²⁵), which says that [if (1) holds] under random changes of $V(0)$, sets of measure zero do not matter. It is precisely pathological behavior on sets of measure zero which are responsible for singular continuous spectra in those almost-periodic models where they occur.¹⁷ The difference between random and almost-periodic models is the decoupling of infinity [which is responsible for (1)] and $V(0)$.²⁸

Finally, we describe some aspects of the proof of theorem 1. Like so much in the one-dimensional

theory, it depends on an analysis of the transfer matrix, $\Phi(n)$, which takes data for solutions of the time-independent Schrödinger equation at 0 to data at n ; i.e.,

$$\Phi(n) = A(n) \cdots A(1); \quad A(j) = \begin{pmatrix} E - V(j) & -1 \\ 1 & 0 \end{pmatrix}. \quad (6)$$

Since the $V(j)$ are random variables, $\Phi(n)$ is a random matrix lying in the group $SL(2, R)$ (for E real). The key technical input for the proof of theorem 1 is that for any k , one can find n so that $\Phi(n)$ has a distribution of the form $G_n(A, E) dA$, where G is C^k in A and E and dA is Haar measure on $SL(2, R)$. This is proven by showing that $G_3(A, E)$ ²⁹ has a fractional derivative in A ,³⁰ noting that G_{3n} is the n -fold $SL(2, R)$ convolution of G_3 , and that repeated convolutions of functions with a fractional derivative are smoother and smoother.³¹

Next, one uses the basic fact noted already by Schmidt³² that one should look at the distribution $d\nu_E(x)$ on $x = u(1)/u(0)$ left invariant by applying an independent random transfer matrix to $(u(1), u(0))$ for³² $k(E) = \int_0^\infty d\nu_E(x)$; so smoothness of $d\nu$ in E implies smoothness of k . Moreover, for each n , ν_E is an eigenfunction of a compact operator built out of G_n . Since the corresponding eigenvalue is simple by results of Furstenberg,³³ the general theory of eigenvalue perturbation theory³⁴ implies that ν_E is at least C^k . Since k is arbitrary, ν_E is C^∞ .

After this paper was submitted, we received two papers from Delyon, Levy, and Souillard,^{35,36} who, also motivated by Kotani,²⁵ discuss localization via a procedure related to, but distinct from, ours in theorems 2 and 3.

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¹B. Simon and M. Taylor, "Harmonic analysis on $SL(2, R)$ and smoothness of the density of states in the one-dimensional Anderson model" (to be published).

²B. Simon and T. Wolff, "Singular continuous spectrum under rank one perturbations and localization for random Hamiltonians" (to be published); B. Simon, "Localization in general one-dimensional random systems" (to be published).

³P. Anderson, Phys. Rev. **109**, 1492 (1958).

⁴In Ref. 1, only a weaker condition on F is needed, namely, that F has compact support and lies in a Sobolev space L_α^1 with $\alpha > 0$. The condition stated in the text is technically simpler to state, implies the L_α^1 condition, and holds for any piecewise C^1 function such as $F(x) = 1$ (0) if $a < x < b$ ($x \leq a$ or $x \geq b$).

⁵L. Pastur, Commun. Math. Phys. **75**, 179 (1980); W. Craig and B. Simon, Commun. Math. Phys. **90**, 207 (1983); F. Delyon and B. Souillard, Commun. Math. Phys. **89**, 415 (1983); *Probability Measures on Groups VII*, edited by E. LePage, Springer Lecture Notes in Mathematics, Vol. 1064 (Springer, Berlin, 1984), p. 309.

⁶F. Wegner, Z. Phys. B **44**, 9 (1981); S. Edwards and D. Thouless, J. Phys. C **4**, 453 (1971); F. Constantinescu, J. Fröhlich, and T. Spencer, to be published.

⁷B. Halperin, Adv. Chem. Phys. **31**, 123-177 (1967).

⁸Any measure $d\kappa$ on $(-\infty, \infty)$ has a unique decomposition $d\kappa = d\kappa^{pp} + d\kappa^{ac} + d\kappa^{sc}$, where the pure point piece is a countable sum of point masses, the absolutely continuous part has the form $g(E)dE$, and the singular continuous part is a measure like the Cantor measure which has no pure points, but lives on a set of Lebesgue measure zero. For spectral measures, the corresponding states are sometimes called "localized," "extended," and "exotic" in the physics literature.

⁹This implies the weak form of localization that for any ϵ there is an R with $\sup_r \sum_{|n| \geq R} |(e^{-iHt} \delta_0)(n)|^2 < \epsilon$ but does not imply boundedness in time of moments like $\sum_n n^2 |(e^{-iHt} \delta_0)(n)|^2$.

¹⁰That is, $d\kappa^{ac}(E) = g(E)dE$ with $g(E)$ almost everywhere nonzero.

¹¹For definiteness, we define the localization length of an eigenfunction ψ to be

$$\lim_{N \rightarrow \infty} \left[(-2N)^{-1} \ln \left[\sum_{|n| \geq N} |\psi(n)|^2 \right] \right]^{-1}.$$

¹²K. Ishii, Prog. Theor. Phys. Suppl. **53**, 77 (1973).

¹³Ishii proved (2) for the half-line problem, but his argument works to control the full-line Green's function; see Ref. 2. We note that, as in Ishii's analysis, (1) implies that

$$\text{Im}[G(0, 0; E_0 + i\epsilon)] = \epsilon \sum_n |G(0, n; E_0 + i\epsilon)|^2$$

goes to zero as $\epsilon \downarrow 0$. This can be used to replace the arguments of F. Martinelli and E. Scoppola, to be published, in the higher-dimensional case.

¹⁴H. Kunz and B. Souillard, Commun. Math. Phys. **78**, 201 (1980); F. Delyon, H. Kunz, and B. Souillard, J. Phys. A **16**, 25 (1983).

¹⁵Ishii's arguments for (2) only depend on positive Lyapunov exponents, and so they imply that (2) holds in the case studied by J. Avron and B. Simon, Duke Math. J. **50**, 369 (1983).

¹⁶Avron and Simon, Ref. 15.

¹⁷W. Craig and B. Simon, Duke Math. J. **50**, 551 (1983).

¹⁸Such results in a related model were first found by R. Carmona, Duke Math. J. **49**, 191 (1982).

¹⁹J. Fröhlich and T. Spencer, Commun. Math. Phys. **88**, 151 (1983).

²⁰Martinelli and Scoppola, Ref. 13.

²¹J. Fröhlich, F. Martinelli, E. Scoppola, and T. Spencer, to

be published; new estimates beyond those in Ref. 19 are required.

²²I. Goldsheid, talk presented at the Conference on Information Theory, Tbilisi, U.S.S.R., September 1984 (unpublished).

²³N. Aronszajn, *Am. J. Math.* **79**, 597 (1957).

²⁴W. Donoghue, *Commun. Pure Appl. Math.* **18**, 559 (1965).

²⁵S. Kotani, in Proceedings of the American Mathematics Society Conference on Random Matrices and Their Applications, Brunswick, Maine, 17–23 June 1984 (to be published).

²⁶See, e.g., M. Reed and B. Simon, *Methods of Modern Mathematical Physics* (Academic, New York, 1972), Vol. I.

²⁷S. Saks, *Theory of the Integral* (Dover, New York, 1964); Y. Katznelson, *An Introduction to Harmonic Analysis* (Dover, New York, 1976).

²⁸This remark can be used to prove a point spectrum in many one-dimensional models with nonindependent $V(n)$;

see Ref. 1.

²⁹ $\Phi(1)$ has a distribution concentrated on a curve in $SL(2, \mathbb{R})$ so that $\Phi(3)$ is the first $\Phi(n)$ that can have a distribution of the form $G(A)dA$ since $SL(2, \mathbb{R})$ is three dimensional. This is why G_3 rather than G_2 appears.

³⁰Even if $d\kappa/dE$ is C^∞ , $G_3(A)$ has an non- L^1 gradient and so the theory of fractional derivatives is essential.

³¹This only explains smoothness in A ; an additional argument (Ref. 1) is needed to get smoothness in E .

³²H. Schmidt, *Phys. Rev.* **105**, 425 (1957).

³³H. Furstenberg, *Trans. Am. Math. Soc.* **108**, 377 (1963).

³⁴T. Kato, *Perturbation Theory for Linear Operators* (Springer, New York, 1966).

³⁵F. Delyon, Y. Levy, and B. Souillard, "Anderson localization for multidimensional systems at large disorder or large energy" (to be published).

³⁶F. Delyon, Y. Levy, and B. Souillard, "An approach 'à la Borland' to multidimensional localization" (to be published).