

18.103, SPRING 2004

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Tuesday and Thursdays, 1-2:30 in Room 4-102. The book is Adams and Guillemin [1]

Lecture 1: February 3.

Read A-G, Chapter 1, pages 1-4, 24-26.

Problems, A-G, Chapter 1, p.11, Problems 1, 2; p.39 Problems 1 (first part), 2 Here are the problems

1.1.1. Prove that the set \mathcal{B} of Bernoulli sequences is uncountable by the Cantor diagonal argument.

1.1.2. a) Let $\omega \in I = (0, 1]$. Show that ω can be written in the form $\sum_{i=1}^{\infty} \frac{a_i}{2^i}$ with $a_i = 0$ or 1 . Show that this expansion is unique when we restrict to nonterminating series.

b) Show that for any (positive) integer k , $\omega \in I$ can be written in the form $\sum_{i=1}^{\infty} \frac{a_i}{k^i}$ where $a_i = 0, 1, \dots, k-1$. Show that the expansion is unique when we restrict to nonterminating series.

1.3.1. Let X be an uncountable set. Let \mathcal{R} be the collection of all finite subsets of X . Given $A \in \mathcal{R}$, let $\mu(A)$ be the number of elements in A . Show that \mathcal{R} is a ring and that μ is a measure on \mathcal{R} .

1.3.2. Let X be an infinite set and let \mathcal{R} be the collection of sets consisting of the finite sets and the complements, in X , of finite sets. Let μ be the function $\mu(A) = 0$ if A is finite, $\mu(A) = 1$ if the complement of A is finite. Is μ a measure?

Lecture 2: February 5.

Read A-G, pp 27-31, pp 4-9

Problems: p.39, nos. 3, 4; p.12 nos. 4,5. Here are the problems

1.3.3. a) Let X be an infinite set and \mathcal{R} the collection of all countable subsets of X . Is \mathcal{R} a ring?

b) Let μ be a measure on \mathcal{R} . Show that there exists a function $f : X \rightarrow [0, \infty)$ such that

$$(1) \quad \mu(A) = \sum_{x \in A} f(x) \quad \forall A \in \mathcal{R}.$$

c) Show that the function f in part b) has the following two properties.

(1) the set $\{x \in X; f(x) \neq 0\}$ is countable and (2) $\sum_{x \in X} f(x) < \infty$.

d) Show that if f has the properties in part c) the formula (1) defines a measure on \mathcal{R} .

1.3.4. Let X be the real line and $\mathcal{R} = \mathcal{R}_{\text{Leb}}$. Given $A \in \mathcal{R}$ let $\mu(A) = 1$ if, for some positive ϵ , A contains the interval $(0, \epsilon)$ and otherwise $\mu(A) = 0$. Show that μ is an additive set function but is *not* countably additive.

1.1.4. Show that

$$\int_0^1 R_{\gamma_1} R_{\gamma_2} \cdots R_{\gamma_n} dx = 0 \text{ or } 1$$

for any sequence $\gamma_1 \leq \gamma_2 \leq \cdots \leq \gamma_n$. When is it 1?

1.1.5. Define the Rademacher functions on the whole real line by requiring them to be periodic of period one – so setting $R_k(x+1) = R_k(x)$. With this definition show that $R_{k+1}(x) = R_k(2x)$ and by induction, $R_k(x) = R_1(2^{k-1}x)$.

Countable additivity and subadditivity

Outer measure

Sets of measure zero

Measurable sets

Rademacher functions

Lecture 3: February 10.

(Special) Chebyshev's inequality

Weak law of large numbers

Lecture 4: February 12.

Homework due Friday 20th: Chapter 1, p. 41, nos 11, 13; p. 42 nos 19, 20.

A-G, pp 33-39.

The σ -ring of measurable sets.

The law of large numbers.

Lecture 5: February 19 (Monday schedule on February 17).

A-G pp 53-58.

Measurable functions

Lecture 6: February 24. The integral

Countable additivity: Theorem 11.

Lecture 7: February 26. Linearity of the integral.

Proof of Theorem 11 finished.

Monotone convergence.

Additivity for positive functions.

The space $L(\mu, \mathcal{F})$ of integrable functions (on X with respect to the measure μ on the σ -ring \mathcal{F}).

Linearity of the integral on integrable functions.

Lecture 8: March 2.

A-G pages 77-80.

Fatou's lemma (for a sequence of non-negative measurable functions)

$$\int_E \liminf_n f_n d\mu \leq \int_E \liminf_n d\mu.$$

Lebesgue's Dominated convergence theorem. Pointwise convergence of a sequence of measurable functions with absolute value bounded by a fixed

integrable function implies convergence of the integrals to the integral of the limit.

The Banach space $L^1(X, \mu)$.

Lecture 9: March 4.

The Banach space $L^1(X, \mu)$ continued.

A Hilbert space, $L^2(X, \mu)$.

Lecture 10: March 9.

Lecture 11: March 11.

Lecture 12: March 16.

Lecture 13: March 18. First in-class test. Covers all material in Lectures 1-10.

Lecture 14: March 30.

Lecture 15: April 1. (I will be away.)

Lecture 16: April 6.

Lecture 17: April 8.

Lecture 18: April 13.

Lecture 19: April 15.

Lecture 20: April 22 (April 20 is a holiday).

Lecture 21: April 27. Second in-class test.

Lecture 22: April 29.

Lecture 23: May 4.

Lecture 24: May 6.

Lecture 25: May 11.

Lecture 26: May 13.

1. OFFICE HOURS

I will try to have an office hour for 18.103 each Wednesday at 3PM. If you are at all desperate to talk, try me after class or send me e-mail to arrange an appointment.

Your intrepid grader, Vito Iacovino, will have an office hour at 6PM on Wednesdays in his office (which is in the basement of the Math Department, down the stairs near 2-190).

2. LECTURE NOTES

2.1. For Lectures 11-14.

I Our discussion of Hilbert space is relevant because of:-

Theorem 1. *For any measure space (X, \mathcal{F}, μ) (countably additive measure on a σ -ring \mathcal{F} on X) the space $L^2(X, \mu)$ is a Hilbert space.*

Proof. We have shown that $L^2(X, \mu)$ is a linear space and that the inner product $\langle f, g \rangle = \int_X f \bar{g} d\mu$ on $L^2(X, \mu)$ is well-defined and satisfies the needed properties. Thus $L^2(X, \mu)$ is a pre-Hilbert space and we only need to show that it is complete for the norm $\|f\|_{L^2} = (\int_X |f|^2 d\mu)^{\frac{1}{2}}$.

We will use the completeness of $L^1(X, \mu)$ proved last week. Let $\{f_n\}$ be a Cauchy sequence in $L^2(X, \mu)$. Consider first the set Z where at least one of the f_n 's is non-zero. This is a countable union of measurable sets,

hence measurable. It is in fact a countable union of sets of finite measure, for instance

$$(1) \quad Z = \bigcup_{n,m} \{|f_n| \geq 1/m\}, \quad \mu\{|f_n| \geq 1/m\} \leq m^2 \int_X |f_n|^2 d\mu.$$

Let us set $Z_k = \bigcup_{n+m \leq k} \{|f_n| \geq 1/m\}$ and so obtain the union as $Z = \bigcup_k Z_k$ where each Z_k has finite measure and $Z_{k+1} \supset Z_k$. Then $g \in L^2(X, \mu)$ implies, by the Cauchy-Schwarz inequality

$$(2) \quad 2 \int_{Z_k} |g| d\mu \leq \mu(Z_k)^{\frac{1}{2}} \|g\|_{L^2}.$$

Applying this to $f_n - f_m$ we see that $\{f_n|_{Z_k}\}$ is Cauchy in $L^1(Z_k, \mu)$ for each fixed k . So, by completeness of L^1 , it converges to $g_k \in L^1(Z_k, \mu)$. Recall that we may change the values of the f_n on a set of measure zero in Z_k so that on a subsequence $f_{n_k(l)}(x) \rightarrow g(x)$ pointwise on Z_k . Passing to the diagonal subsequence in k , the uniqueness of the limit (modulo values on sets of measure zero) means that there is one function g defined on the whole of Z such that $f_{n(l)} \rightarrow g$ almost everywhere. Taking $g = 0$ on $X \setminus Z$ and again changing values (to 0) on a set of measure zero if necessary we can arrange that $f_{n(l)}(x) \rightarrow g(x)$ pointwise on X .

Now that we have our putative limit we can use Fatou's Lemma. Applied to the sequence of measurable, non-negative, functions $|f_m - f_n|^2$ with n fixed and $m \rightarrow \infty$ it states that

$$(3) \quad \int_X \liminf_m |f_n(x) - f_m(x)|^2 d\mu \leq \liminf_m \int_X |f_m - f_n|^2 d\mu.$$

Given $\epsilon > 0$ the fact that $\{f_n\}$ is Cauchy in $L^2(X, \mu)$ allows us to choose N such that for $n, m \geq N$ $\|f_n - f_m\|_{L^2} < \epsilon$. This implies that, for $n \geq N$ the right side of (3) is less than or equal to ϵ^2 . On the other hand the integrand on the left side is pointwise convergent so

$$(4) \quad \int_X |f_n(x) - g|^2 d\mu \leq \epsilon^2 \quad \forall n \geq N.$$

This shows that $f_n \rightarrow g$ in $L^2(X, \mu)$ (and in particular that $g \in L^2(X, \mu)$). \square

II Now, suppose we are in a general Hilbert space H . Suppose that $\{\phi_i\}_{i \in I}$ is a countable – either finite or countably infinite – orthonormal set. That is, each element $\phi_i \in H$ has norm one $\|\phi_i\|^2 = \langle \phi_i, \phi_i \rangle = 1$ and they are pairwise orthogonal so $\langle \phi_i, \phi_j \rangle = 0$ if $i \neq j$. Then for each element $f \in H$ we can consider the constants

$$(5) \quad c_n(f) = \langle f, \phi_n \rangle, \quad n \in I.$$

Proposition 1. *For any orthonormal set and any $f \in H$ the series*

$$\sum_{i \in I} c_n(f) \phi_n$$

converges in H and

$$(6) \quad \sum_{i \in I} |c_i(f)|^2 = \left\| \sum_{i \in I} c_i(f) \phi_n \right\|^2 \leq \|f\|^2.$$

[This is Bessel's inequality.]

Proof. Take an ordering of I so that we can replace it by $\{1, \dots\}$ and for any finite N consider the finite sum

$$(7) \quad S_N(f) = \sum_{n=1}^N c_n(f) \phi_n \in H.$$

Expanding out the inner product using the sesquilinearity gives

$$(8) \quad \|S_N(f)\|^2 = \langle S_N(f), S_N(f) \rangle = \sum_{i,j=1}^N c_i(f) \overline{c_j(f)} \langle \phi_i, \phi_j \rangle = \sum_{i=1}^N |c_i(f)|^2$$

using the orthonormality of the ϕ_i 's. Doing the same thing for $S_M(f) - S_N(f)$ where $M \geq N$ we find

$$(9) \quad \|S_m(f) - S_M(f)\|^2 = \sum_{n=N}^M |c_n(f)|^2.$$

On the other hand if we write $f = (f - S_N(f)) + S_N(f)$ we see that $\langle f - S_N(f), \phi_n \rangle = 0$ for $n \leq N$. Since $S_N(f)$ is a sum of these ϕ_n 's, $\langle f - S_N(f), S_N(f) \rangle = 0$. This in turn means that

$$(10) \quad \|f\|^2 = \langle (f - S_N(f)) + S_N(f), (f - S_N(f)) + S_N(f) \rangle \\ = \|f - S_N(f)\|^2 + \|S_N(f)\|^2$$

since the cross-terms vanish in the linear expansion. Thus

$$(11) \quad \|S_N(f)\|^2 \leq \|f\|^2 \quad \forall N.$$

Going back to (8) we conclude from (11) that the series $\sum_{n=1}^N |c_n|^2$ has an upper bound independent of N . When I is infinite this means it converges (in \mathbb{R}). From (9) it follows that the sequence $\{S_N(f)\}$ is Cauchy in H . Since H is complete, being a Hilbert space, it must converge and the limit must satisfy (6). \square

See if you can show that the limit $\sum_{i \in I} c_i(f) \phi_i$ is independent of the order chosen for I .

III A (countable) orthonormal set is said to be complete if $S(f) = f$ for all $f \in H$. Notice that $f - S(f)$ is the limit of $f - S_N(f)$ as $N \rightarrow \infty$ and $\langle f - S_N(f), \phi_j \rangle = 0$ whenever $j \geq N$. Taking the limit we see that

$$(12) \quad \langle f - S(f), \phi_j \rangle = \lim_{N \rightarrow \infty} \langle f - S_N(f), \phi_j \rangle = 0$$

for all j , where we have used the fact that if $v_j \rightarrow v$ in H then $\langle v_j, w \rangle \rightarrow \langle v, w \rangle$ for each $w \in H$ - this follows from Schwarz inequality since

$$(13) \quad |\langle v, w \rangle - \langle v_j, w \rangle| \leq \|v - v_j\| \|w\|.$$

Thus another way of stating the completeness of the orthonormal set is that

$$w \in H, \langle w, \phi_j \rangle = 0 \quad \forall j \implies w = 0.$$

IV Existence of complete orthonormal bases.

Theorem 2. *Any separable Hilbert space has a (countable) complete orthonormal basis.*

Proof. Recall that a metric space is separable if it has a countable dense subset. So we can suppose there is a countable set $E \subset H$ with $\bar{E} = H$. Let $E = \{e_1, e_2, \dots\}$ be an enumeration of E . We extract a complete orthonormal basis from E by applying the Gram-Schmid procedure. First consider e_1 . If it is zero, pass to e_2 . If it is non-zero, set $\phi_1 = e_1/\|e_1\|$ which gives an orthonormal set with one element, now pass to e_2 . Proceeding by induction, suppose at stage n we have an orthonormal set $\{\phi_1, \dots, \phi_N\}$ with $N \leq n$ elements such that each of the e_j , $j \leq n$, is linearly dependent on these N elements. Now consider e_{n+1} . If it is dependent on the ϕ_j for $j \leq N$, pass on to e_{n+2} . If not then

$$(14) \quad \phi_{N+1} = g/\|g\|, \quad g = e_n - \sum_{j=1}^N \langle g, \phi_j \rangle \phi_j$$

is well-defined, such that $\{\phi_1, \dots, \phi_N\} \cup \{\phi_{N+1}\}$ is orthonormal and such that the e_j for $j \leq N+1$ are dependent on this new orthonormal set.

Thus we can proceed by induction to define an orthonormal set, which will either be finite or countable (depending on H). In either case it is complete. To see this, suppose there is some element $f \in H$ orthogonal to all the ϕ_i . By the density of E , for any $\epsilon > 0$ there exists $e \in E$ such that $\|f - e\| < \epsilon$. However, e is in the (finite) span of the ϕ_i , so $\langle f, e \rangle = 0$. This however implies, by Pythagorus' theorem, that $\epsilon^2 \geq \|f - e\|^2 = \|f\|^2 + \|e\|^2$. Thus $\|f\| \leq \epsilon$ so in fact $\|f\| = 0$ and hence $f = 0$, proving the completeness. \square

V The basic result we will prove on Fourier series is that for the special case of $L^2([-\pi, \pi])$, computed with respect to Lebesgue measure, the exponentials

$$(15) \quad \phi_n(x) = \frac{1}{\sqrt{2\pi}} \exp(inx), \quad n \in \mathbb{Z},$$
 form a complete orthonormal set.

First we want to check that it is indeed an orthonormal set. Since $|\exp(inx)| = 1$ the norm is easy enough to compute

$$(16) \quad \|\phi_n\|_{L^2}^2 = \int_{[-\pi, \pi]} \frac{1}{2\pi} = 1$$

since we do know how to integrate constants. The orthogonality would seem almost as easy,

$$(17) \quad \langle \phi_n, \phi_k \rangle = \frac{1}{2\pi} \int_{[-\pi, \pi]} e^{i(n-k)x} dx = \frac{1}{2\pi} \frac{e^{i(n-k)x}}{i(n-k)} \Big|_{-\pi}^{\pi} = 0.$$

However, this is proof by abuse of notation since here we are using the Riemann integral (and Fundamental Theorem of Calculus) and we are supposed to be computing the Lebesgue integral.

So I need to go back and check some version of their equality. The following will do for present purposes.

Proposition 2. *If g is a continuous function on a finite interval $[a, b] \subset \mathbb{R}$ then its Riemann and Lebesgue integrals are equal.*

Proof. We can split g into real and imaginary parts if it is complex valued and the result then follows from the real case; so assume that g is real. Since we do know how to integrate constants (being simple functions) we can add to g the constant $-\inf_{[a,b]} g$ (or something larger) and we can then assume that $g \geq 0$. Now, the Riemann integral is defined as the common value of the upper and lower integrals (look this up in Rudin [2], I am not going to remind you of all of it.) One result for continuous functions (in fact general Riemann integrable functions) is that given $\epsilon > 0$ there exist a partition of $[a, b]$, \mathcal{P} , such that the difference of lower and upper partial sums satisfies

$$(18) \quad U(g, \mathcal{P}) - L(g, \mathcal{P}) < \epsilon.$$

Notice here that the lower sum is actually $I_{[a,b]}(s)$ for a simple function which is smaller than g . So, directly from the definition of the integral we know that

$$(19) \quad \int_a^b g dx = \sup_{\mathcal{P}} L(g, \mathcal{P}) \leq \int_{[a,b]} g dx$$

where the integral on the left is Riemann's and on the right is Lebesgue's. On the other hand if we simply divide $[a, b]$ into N equal parts and take the simple function with value $\inf g$ on each interval we get a sequence of simple functions increasing to g (in fact uniformly on $[a, b]$.) The Riemann lower partial sum for this partition is $I(g_N)$, bounded above by the Riemann integral, and by the monotone convergence theorem this sequence converges to the Lebesgue integral. This gives the opposite inequality to (19) so the two integrals are equal. \square

This argument only needs slight modification to show that every Riemann integrable function on $[a, b]$ is Lebesgue integrable and that the integrals are equal; it is done in Adams and Guillemin.

Thus we know that the Fourier functions $\phi_n(x)$ do indeed form a countable orthonormal set for $L^2([-\pi, \pi])$. We still need to know that it is complete. This involves some more work.

VI To prove the completeness of the Fourier functions $\{\phi_n\}$ we need to show that any function $f \in \mathcal{L}^2([-\pi, \pi])$ which satisfies

$$(20) \quad c_n(f) = \frac{1}{\sqrt{2\pi}} \int_{[-\pi, \pi]} f(x) e^{-inx} dx = 0 \quad \forall n \in \mathbb{Z},$$

itself vanishes almost everywhere, so is zero in $L^2([-\pi, \pi])$. In fact we will show something a little stronger,

Proposition 3. *If $f \in \mathcal{L}^1([-\pi, \pi])$ satisfies (20) then $f = 0$ outside a set of measure zero in $[-\pi, \pi]$.*

However this will take some work.

VII First we make the following observation directly from our construction of the integral.

Lemma 1. *If $f \in \mathcal{L}^1(X, \mu)$ and $A_j \subset X$ is a sequence of measurable sets and $A \subset X$ with $\mu(S(A, A_j)) \rightarrow 0$ then*

$$(21) \quad \int_{A_j} f d\mu \longrightarrow \int_A f d\mu.$$

Proof. This is really just a reminder of what we have done earlier. We might as well assume that f is non-negative, since we can work with the real and imaginary parts, and then with their positive and negative parts, separately. It suffices to show that every subsequence of the real sequence in (21) has a convergent subsequence with limit the integral over A . Since we are not assuming anything much about the A_j 's it is enough to show that there is a subsequence converging to the integral over A and then apply the argument to any subsequence. Since the measure of the symmetric difference $\mu(S(A_j, A)) \rightarrow 0$ we can pass to a subsequence (which we then renumber) so that

$$(22) \quad \mu(S(A_j, A)) = \mu(A \setminus A_j) + \mu(A_j \setminus A) \leq 2^{-j}.$$

Since $\int_A g d\mu - \int_{A_j} g d\mu = \int_{A \setminus A_j} g d\mu - \int_{A_j \setminus A} g d\mu$ it suffices to show that these sequences tend to zero. So, set $B_j = A \setminus A_j$; from (22) it follows that $F = \bigcup_{j \geq 1} B_j$ has finite measure and $F_N = \bigcup_{j \geq N} B_j$ has measure tending to zero, with $B_N \subset F_N$. Thus it suffices to show that $\int_{F_N} g d\mu \rightarrow 0$ since this sequence dominates $\int_{B_j} g d\mu$ by monotonicity. Finally then we write $F = \bigcup_{j \geq 1} G_j$, $G_j = F_j \setminus F_{j+1}$ which is a decomposition into a countable collection of disjoint open sets. By the countable additivity of the integral

$$(23) \quad \int_F g d\mu = \sum_{j=1}^{\infty} \int_{G_j} g d\mu.$$

Thus the series of non-negative terms on the right converges which implies that the series 'of remainders'

$$(24) \quad \int_{F_N} g d\mu = \sum_{j=N}^{\infty} \int_{G_j} g d\mu \rightarrow 0 \text{ as } N \rightarrow \infty.$$

This shows that the $\int_{A \setminus A_j} g d\mu \rightarrow 0$ as $j \rightarrow \infty$; the other half of (22) can be handled in the same way, so we have a subsequence with the correct convergence and hence have proved the Lemma. \square

VIII So, let's apply Lemma 1 directly as follows.

Lemma 2. *If $f \in L^1([a, b])$ for a finite interval $[a, b] \subset \mathbb{R}$ then $g(s) = \int_{[a, s]} f dx$ is a continuous function on $[a, b]$. Furthermore, if $f_i \rightarrow f$ in $L^1([a, b])$ then the corresponding functions $g_i \rightarrow g$ uniformly on $[a, b]$.*

Proof. By Lemma 1, $\int_{[a, s+t]} g dx \longrightarrow \int_{[a, s]} g dx$ as $t \rightarrow 0$. \square

IX Now, go back and consider $f \in L^1([-\pi, \pi])$ which satisfies (20). We will show that there is a continuous function (not obviously zero) which also satisfies (20). Namely consider

$$(25) \quad g(s) = \int_{[-\pi, s]} f dx - C, \quad C \text{ chosen so } \int_{[-\pi, \pi]} g dx = 0.$$

The constant term here is added so that $c_0(g) = 0$. We have to work harder to show that the other $c_n(g) = 0$. To do so we use the integration by parts identity

$$(26) \quad \int_{[a,b]} h_1(s) \left(\int_{[a,s]} h_2(x) dx \right) ds = \int_{[a,b]} \left(\int_{[x,b]} h_1(s) ds \right) h_2(x) dx.$$

Proposition 4. *The identity (26) holds if $h_i \in L^1([a, b])$, $i = 1, 2$.*

Before worrying about the proof of this, let us apply it to $h_2 = f$ and $h_1 = \exp(-inx)$ for some $0 \neq n \in \mathbb{Z}$. On the left in (26) the inner integral is our definition of g in (25) except for the missing constant. On the right in (26) we compute, using equality of Riemann and Lebesgue integrals for continuous functions, finding

$$(27) \quad \int_{[b,x]} h_1(s) ds = -\frac{1}{in} (\exp(-inb) - \exp(-inx)).$$

This is a linear combination of our Fourier exponentials, so

$$(28) \quad \int_{[a,b]} \left(\int_{[x,b]} \exp(-ins) ds \right) f(x) dx = 0, \quad 0 \neq n \in \mathbb{Z} \text{ if } f \text{ satisfies (20).}$$

Since we already have arranged that $c_0(g) = 0$ and the constant does not change the $c_n(g)$ for $n \neq 0$, once we prove Proposition 4 we will know that

$$(29) \quad g \text{ given by (25) satisfies (20) if } f \text{ does so.}$$

X Proof of Proposition 4.

We have to prove (26). Again this is a return to basics. Notice first that both sides of (26) do make sense if $h_i \in L^1([a, b])$ since then the integrated functions

$$(30) \quad g_2(s) = \int_{[a,s]} h_2(x) dx \text{ and } g_1(x) = \int_{[x,b]} h_1(s) ds$$

are both continuous. This means that the products $h_1(s)g_2(s)$ and $h_2(x)g_1(x)$ are both integrable (why exactly).

Now, (26) is separately linear in h_1 and h_2 . So as usual we can assume these functions are positive, by first replacing the functions by their real and imaginary parts, and then their positive and negative parts. On the left the we can take a sequence of non-negative simple functions approaching h_1 from below and we get convergence in $L^1([a, b])$. Similarly on the right the integrated functions approach g_1 uniformly so the integrals converge. The same argument works for h_2 , so it suffices to prove (26) for simple functions and hence, again using linearity, for the characteristic functions of two measurable sets, $h_i = \chi_{A_i}$, $A_i \subset [a, b]$ measurable. Recalling what it means to be measurable, we can approximate say A_1 , in measure, by a sequence of elementary sets, each a finite union of disjoint intervals. Using Lemma 1 (several times) the resulting integrals converge. Applying the same argument to A_2 it suffices to prove (26) when h_1 and h_2 are the characteristic functions of intervals. In this case our identity has become rather trivial, except that we have to worry about the various cases. So, suppose that h_1 is the characteristic function of $[a_1, b_1]$, a subinterval of

$[a, b]$ (open, half-open or closed does not matter of course). We are trying to prove

$$(31) \quad \int_{[a_1, b_1]} \left(\int_{[a, s]} h_2(x) dx \right) ds \stackrel{?}{=} \int_{[a, b]} \left(\int_{[x, b] \cap [a_1, b_1]} ds \right) h_2(x) dx.$$

Using the constraints on the domains (these are now Riemann integrals anyway) this is equivalent to

$$(32) \quad \int_{[a_1, b_1]} \left(\int_{[a_1, s]} \chi_{[c, d]}(x) dx \right) ds \stackrel{?}{=} \int_{[a_1, b_1]} \left(\int_{[x, b_1]} ds \right) \chi_{[c, d]}(x) dx$$

We can replace a_1 and c both by $\max(a_1, c)$ from the support properties and similarly we can replace b_1 and d by $\min(b_1, d)$. Calling the new end points a and b again we are down to

$$(33) \quad \int_{[a, b]} \left(\int_{[a, s]} dx \right) ds \stackrel{?}{=} \int_{[a, b]} \left(\int_{[x, b]} ds \right) dx$$

which is just a very special case of Riemannian integration by parts. Thus (26), and hence Proposition 4, is proved in general.

XI Okay, so now we know that if $f \in L^1([-\pi, \pi])$ has all the $c_n(f) = 0$ then g in (25) is continuous and satisfies the same thing. Of course we can continue this, and integrate again, replacing g by f and get a new function which satisfies

$$(34) \quad h(s) = \int_{[-\pi, s]} g dx - C', \quad h' = g, \quad c_n(h) = 0 \quad \forall n \in \mathbb{Z}.$$

Thus h is continuously differentiable and still has $c_n(h) = 0$. We will show that this implies $h \equiv 0$ by proving a convergence result for Fourier series, this is [1] Theorem 2 on p.140, more or less.

XII

Theorem 3. *If h is differentiable on $[-\pi, \pi]$ (so only from the right at $-\pi$ and from the left at π) with continuous derivative then*

$$(35) \quad S_n(h)(x) \longrightarrow h(x) \quad \forall x \in [-\pi, \pi].$$

Proof. Pick a point $x_0 \in [-\pi, \pi]$ and consider the partial sum of the series we are interested in

$$(36) \quad S_N(h)(x_0) = \frac{1}{\sqrt{2\pi}} \sum_{k=-N}^N c_k(h) e^{ikx_0} = \frac{1}{2\pi} \int_{-\pi}^{\pi} h(x) \left(\sum_{k=-N}^N e^{ik(x_0-x)} \right) dx.$$

Here we have just inserted the definition of the $c_n(h)$. So consider the function

$$(37) \quad D_N(x) = \frac{1}{2\pi} \sum_{k=-N}^N e^{ikx} = \frac{1}{2\pi} \frac{e^{i(N+1)x} - e^{-iNx}}{e^{ix} - 1}.$$

To see this, multiply the sum defining $D_N(x)$ by e^{iNx} and observe that it becomes $\sum_{i=0}^{2N} T^k = (T^{2N+1} - 1)/(T - 1)$ where $T = e^{ix}$. Also integrating $D_N(x)$ term by term we find that

$$(38) \quad \int_{-\pi}^{\pi} D_N(x) dx = 1, \quad D_N(x + 2\pi n) = D_N(x), \quad x \in \mathbb{R}, \quad n \in \mathbb{Z}.$$

The reason for looking at $D_N(x)$ is that, from (36),

$$(39) \quad S_N(h)(x_0) = \int_{-\pi}^{\pi} h(x)D_N(x_0 - x)dx.$$

From (38) we can write the expected limit in a similar way

$$(40) \quad h(x_0) = \int_{-\pi}^{\pi} h(x_0)D_N(x_0 - x)dx$$

where we use the second part of (38) to reorganize the integral. Combining these two we get

$$(41) \quad [S_N(h)(x_0) - h(x_0)] = \int_{-\pi}^{\pi} [h(x) - h(x_0)]D_N(x_0 - x)dx.$$

Now, the function $e^{ix} - 1$ in the denominator of $D_N(x)$ in (37) vanishes precisely at $x = 0$ in $[-\pi, \pi]$ and does so simply. That is, $x/(e^{ix} - 1)$ is continuous if we define it to take the value of the derivative $-i$ at $x = 0$. This means that the quotient

$$(42) \quad p(x) = \frac{h(x) - h(x_0)}{e^{i(x-x_0)} - 1} = \frac{h(x) - h(x_0)}{x - x_0} \frac{x - x_0}{e^{i(x-x_0)} - 1}$$

is bounded and continuous on $[-\pi, \pi]$ if we define it correctly at $x = x_0$. In particular this means it defines an element $p \in L^2([-\pi, \pi])$. Then (41) can be written

$$(43) \quad [S_N(h)(x_0) - h(x_0)] = c_{N+1}(p) - c_{-N}(p), \quad p \in L^2([-\pi, \pi]).$$

We already know, by Bessel's inequality, that the series $\sum_{n \in \mathbb{Z}} |c_n(p)|^2$ converges so $c_{N+1}(p), c_{-N}(p) \rightarrow 0$ as $N \rightarrow \infty$. Thus we have proved Theorem 3. \square

I should have commented a little on the case $x_0 = \pm\pi$.

XII So, now we know that $S_N(h)(x_0) \rightarrow H(x_0)$ for each $x_0 \in [-\pi, \pi]$ if h is differentiable. In fact the proof shows a little more than this, so let us record it:

Lemma 3. *If $h \in L^2([\pi, \pi])$ is continuous and differentiable from the right and left at x_0 then $S_N(h)(x_0) \rightarrow H(x_0)$.*

Proof. Just check that this is all we really used. \square

XIII Returning to our efforts to prove Proposition 3 we now know that if f has $c_n(f) = 0$ for all n then g given by (25) must be constant, since its indefinite integral h vanishes. Thus the integrals $\int_{[-\pi, s]} f(x)dx$ are independent of the end point s , so must all vanish (since we know that the limit as $s \downarrow -\pi$ is zero). Thus we have a function $f \in L^1([-\pi, \pi])$ with integral zero. The proof of Proposition 3 is therefore finished by

Lemma 4. *If $f \in L^1([a, b])$ and $\int_a^x f(s)ds = 0$ for all $x \in [a, b]$ then $f = 0$ in $L^1([a, b])$.*

Proof. First, we can take the difference

$$(44) \quad \int_{[c,d]} f(s)ds = \int_a^d f(s)ds - \int_a^c f(s)ds = 0$$

to conclude that

$$(45) \quad \int_A f dx = 0 \quad \forall A \in \mathcal{R}_{\text{Leb}},$$

the ring of subsets of $[a, b]$ consisting of finite unions of disjoint intervals. But then, by Lemma 1, we see that the integral vanishes for all Lebesgue subsets of $[a, b]$, since they can be approximated by such finite unions. Now, take A to be the measurable set on which $f \geq 0$ and then the set on which $f \leq 0$ and we conclude that $f = 0$ almost everywhere in $[a, b]$ and hence $f = 0$ in $L^1([a, b])$. \square

XIV Finally then we have proved Proposition 3 which can be restated as

(46) No two distinct elements of $L^1([-\pi, \pi])$ have the same Fourier coefficients,

or, if you prefer, that the Fourier coefficients determine a function in $L^1([-\pi, \pi])$. For a function in $L^2([-\pi, \pi])$ we have the stronger statement that

$$(47) \quad f(x) = \sum_{n=-\infty}^{\infty} c_n(f) e^{inx} \text{ converges in } L^2([-\pi, \pi])$$

$$\text{with } c_n(f) = \frac{1}{\sqrt{2\pi}} \int_{-\pi}^{\pi} f(x) e^{-inx} dx.$$

3. HOMEWORK

Always due on Fridays by Noon, in 2-108 (there will be a box somewhere near the door for 18.103).

- (1) due February 6: A-G, Chapter 1, p.11, nos. 1, 2; p.12 nos. 4,5; p.39 nos 1 (first part), 2, 3, 4. [Since this is all a bit rushed, if you prefer you can leave p.12 4,5 and p.39 3,4 until next Friday, Feb 13.]
- (2) due February 13: A-G, Chapter 1, p.12 4,5 and p.39 3,4 from last time.
- (3) due February 20: A-G, Chapter 1, p. 41, nos 11, 13; p. 42 nos 19, 20.
- (4) due February 27: I failed to put anything up in time (I blame the disasterous and terminal crash of my laptop).
- (5) due March 5:: A-G, Chapter 2, p. 58, nos 3, 8, 10; p. 69, no 2; p.70 no 9; p.80 no 1.
- (6) due March 12:
- (7) due April 2.
- (8) due April 9.
- (9) due April 16.
- (10) due April 23.
- (11) due May 7.

4. TESTS AND FINAL EXAM

First Test, March 18.

Second Test, April 27.

Final Exam, Exam period.

5. PRACTICE FOR FIRST TEST

I have decided to make the first test open-book, but only copies of Adams and Guillemin are permitted, no papers or other writings.

- (1) Set $X = [0, 1] \times [0, 1]$, the unit square in \mathbb{R}^2 and suppose that $f : X \rightarrow [0, \infty)$ is Lebesgue measurable and such that

$$\mu_{\text{Leb}}\{x \in X; f(x) \geq n\} \leq 2^{-n} \quad \forall n \in \mathbb{N}.$$

Show that $\int_X f dx < \infty$.

- (2) Without using the equality of Riemann and Lebesgue integrals show that

$$\int_{\mathbb{R}} (1 + |x|)^{-1} dx = \infty$$

and

$$\int_{\mathbb{R}} (1 + |x|^2)^{-1} dx < \infty.$$

- (3) Give, with proofs, an example of a non-negative Lebesgue measurable function on $[0, 1]$ which has finite Lebesgue integral but is not bounded.
 (4) Explain why the formula

$$\mu(A) = \int_A \exp(-x^2) dx$$

defines a countably additive finite measure on the σ -ring, \mathcal{M} , of all Lebesgue measurable subsets of \mathbb{R} . Show that completion of \mathcal{M} with respect to this measure is just \mathcal{M} again.

- (5) Let (X, \mathcal{F}, μ) be a measure space and suppose that $s_n : X \rightarrow [0, \infty)$ is an increasing (i.e. pointwise non-decreasing) sequence of simple measurable functions with $\int_X s_n d\mu < 1$. Show that the set consisting of those points $x \in X$ where $s_n(x) \rightarrow \infty$ has μ -measure zero.

6. GRADES

Grades mean what they are supposed to:

- A Essential mastery of the course.
- B Clear facility with most of the material.
- C Marginal understanding of much of the material.
- D Some comprehension.

REFERENCES

- [1] Malcolm Adams and Victor Guillemin, *Measure theory and probability*, Springer Verlag, 1996.
 [2] W. Rudin, *Real and complex analysis*, third edition ed., McGraw-Hill, 1987.

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