INVARIANT MEASURES ON PRODUCT SPACES

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1. Introduction

Let G be a locally compact group (which we always assume to satisfy the second axiom of countability), and let X be a standard Borel space on which G acts as a Borel transformation group (see [7], p. 628). That is, we have a homomorphism of G into the group of Borel automorphisms of X such that if $x \to g \cdot x$ is the automorphism corresponding to $g \in G$, then $(x, g) \to g \cdot x$ is a Borel function from $G \times X$ into X where G is endowed with the σ -field of sets generated by the open sets, and $G \times X$ is given the product σ -field. Further, let μ be a σ -finite measure on X (all measures henceforth will be understood to be σ -finite) which is quasi-invariant under G; that is, for every $g \in G$, $g \cdot \mu$ and μ are equivalent in the sense of mutual absolute continuity. (Here $g \cdot \mu$ is the transform of μ by g defined by $(g \cdot \mu)(\sigma) = \mu(g^{-1} \cdot \sigma)$ for Borel sets σ of X.) One says that μ is ergodic under G if for every Borel σ in G such that $\mu(\sigma\Delta g \cdot \sigma) = 0$ for all $g \in G$, we have $\mu(\sigma) = 0$ or $\mu(X - \sigma) = 0$. It is clear that these two properties of μ depend not on μ , but only on the equivalence class $C(\mu)$ of μ . We shall say, following [8], that $C(\mu)$ is a quasi-orbit of G if μ is quasi-invariant and ergodic. Note that each orbit of G on X carries a unique equivalence class $C(\mu)$ of such measures (see [8], p. 295). One calls these classes transitive quasi-orbits or simply orbits.

We shall say that a measure ν is invariant if $g \cdot \nu = \nu$ for all $g \in G$. In this note we are going to discuss a special case of the following circle of questions: given a quasi-orbit $C(\mu)$ on X, when does it contain an invariant measure ν or more specifically an invariant ν with specified properties? We note that if $\nu \in C(\mu)$ is invariant, it is unique up to multiplication by positive scalars. This is an immediate consequence of ergodicity. Furthermore, any $\lambda \in C(\mu)$ is either atomic (consists of point masses) or nonatomic (no point masses).

The systems $(G, X, C(\mu))$ which we will discuss will be such that G is countable and acts freely on X in the sense that $\{x: g \cdot x = x \text{ for some } g \neq e\}$ is a μ -null set. When these conditions are satisfied, von Neumann in [10] has shown how to construct a certain factor von Neumann algebra associated with $(G, X, C(\mu))$. The type of this factor is determined by the measure theoretic properties of $C(\mu)$ discussed in the previous paragraph (see [10], theorem IX). One can show

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that these kinds of factors arise as the rings generated by certain representations of restricted products of groups. Our results will be applied to this situation in a subsequent paper.

Our attention will be restricted to those systems $(G, X, C(\mu))$ which arise in a very particular fashion. Let A be a countable index set, and let X_{α} , $\alpha \in A$ be countable spaces with the σ -field of all subsets, and let G_{α} be a discrete group acting transitively on X_{α} . Let μ_{α} be a probability measure on X_{α} , quasi-invariant under G_{α} . (This simply means, of course, that each point of X_{α} gets positive mass.) Let X be the full Cartesian product of the X_{α} given the product Borel structure (which is a standard Borel space), and let μ be the infinite product of the measures μ_{α} (see [4]; recall that $\mu_{\alpha}(X_{\alpha}) = 1$). Finally let G be the direct sum of G_{α} ; $G = \{(g_{\alpha}) : g_{\alpha} \in G_{\alpha}, g_{\alpha} = e$ for almost all α }. ('Almost everywhere' with respect to sets of indices α will mean 'except for a finite number'; there should be no confusion with the usual measure theoretic usage of the term.) Then G is countable and acts on X in the obvious way. Moreover it is clear (see below) that $C(\mu)$ is a quasi orbit. Our results concern conditions under which $C(\mu)$ contains an invariant measure.

Von Neumann raised this question in the case when X_{α} is a two-point space, and implicitly states a theorem but gives no details (see [10], p. 95, last sentence of the third paragraph of 4) and compare with ([9], pp. 66–77). Pukansky treated a special case of this, ([11], p. 140, lemma 6), and recently J. D. C. Bures [2] has given a systematic treatment of some aspects of this including a proof of von Neumann's assertion. See also the paper of Araki [1]. Bures' results are unfortunately not strong enough for our purposes. Bures works in the context of the von Neumann algebra associated to the system, whereas we shall work just with the measure theory. We wish to thank R. V. Kadison for calling our attention to this paper of Bures.

2. Statement of the results

Let $(G, X, C(\mu))$ be a triple as discussed above, not necessarily arising from the product construction, with G countable and acting freely where $C(\mu)$ is a quasi-orbit.

DEFINITIONS. (1) We shall say that $C(\mu)$ is type I if it is transitive (or equivalently that $C(\mu)$ consists of atomic measures, or that $C(\mu)$ contains an atomic invariant measure).

- (2) We shall say that $C(\mu)$ is type II_1 if it is nontransitive (that is, nonatomic), and contains a finite invariant measure.
- (3) We shall say that $C(\mu)$ is type II_{∞} if it is nontransitive, and contains an infinite invariant measure.
 - (4) We shall say that $C(\mu)$ is type III if it contains no invariant measure.

We observe that these possibilities are mutually exclusive and exhaustive. One could, of course, formulate these definitions for an arbitrary G, but we shall not need this, and moreover, if G does not act freely, this classification does not

seem to be too appropriate (see [8]). We note that the von Neumann algebra constructed from $(G, X, C(\mu))$ mentioned in section 1 is a factor whose type is exactly that of $C(\mu)$ ([10], theorem IX). Our approach here is to fix G and let $C(\mu)$ vary. However, the recent result of Tulcea [13] concerns the opposite. This says roughly that if $C(\mu)$ is fixed and G is fixed to be the integers, but its action is allowed to vary, then in a suitable categorical sense $C(\mu)$ is almost always type III. Our results indicate the same, namely that the type III case is generic.

Let us now return to the specific case which we shall consider. That is, $X = \prod_{\alpha} X_{\alpha}$, $\mu = \prod_{\alpha} \mu_{\alpha}$ where the X_{α} are countable and the μ_{α} are probability measures on X_{α} giving each point positive mass, and G is the direct sum of countable groups G_{α} with G_{α} transitive and free on X_{α} .

LEMMA 2.1. The class $C(\mu)$ is a quasi-orbit for G.

PROOF. This is essentially known, but we include a proof for completeness. If $\alpha \in A$, then $X = X_{\alpha} \times X_{\alpha}^*$ where X_{α}^* is the product of the X_{β} , $\beta \neq \alpha$. Also, $\mu = \mu_{\alpha} \times \mu_{\alpha}^*$ where μ_{α}^* is the product of the remaining μ_{β} , $\beta \neq \alpha$. Now let us identify G_{α} with the subgroup of G, all of whose coordinates except the α -th are the identity element. Then G_{α} acts on X by its usual action on X_{α} and trivially on X_{α}^* . Clearly μ is quasi-invariant under G_{α} since μ_{α} is. Since G is generated by the G_{α} , μ is quasi-invariant under G. Finally, to see that $C(\mu)$ is ergodic, we simply note that this assertion is the Borel zero-one law of probability theory (see [3], p. 102).

We next want to note that the type of $C(\mu)$ does not depend on any precise information about the G_{α} , since one can characterize G-invariant measures on X independently of the G_{α} .

LEMMA 2.2. Let ν be a Borel measure on X. Then ν is invariant under G if and only if for each α , $\nu = \nu_{\alpha} \times \nu_{\alpha}^{*}$, where ν_{α}^{*} is some Borel measure on X_{α}^{*} and ν_{α} is the measure on X_{α} , giving mass one to each point.

Proof. If our condition is satisfied, then $\nu = \nu_{\alpha} \times \nu_{\alpha}^{*}$, and since ν_{α} is invariant under G_{α} , ν is invariant under G_{α} since G_{α} acts trivially on X_{α}^{*} . Thus again, as G is generated by the G_{α} , ν is G-invariant. Conversely, let ν be G-invariant, and write $X = X_{\alpha} \times X_{\alpha}^{*}$. Let $x \in X_{\alpha}$ and define a measure ν_{α}^{*} on X_{α}^{*} by $\nu_{\alpha}^{*}(\sigma) = \nu(\{x\} \times \sigma)$ for $\sigma \subset X_{\alpha}^{*}$. It is easy to see that ν_{α}^{*} is σ -finite so that it is a measure under our conventions. Since G_{α} is transitive on X_{α} and ν is invariant under G_{α} , one sees that ν_{α}^{*} is independent of the choice of x. As X_{α} is countable, it is clear that $\nu = \nu_{\alpha} \times \nu_{\alpha}^{*}$ on all 'rectangles' in the product. Therefore, these two measures are the same. This completes the proof.

We note indeed that the condition of the lemma is independent of the G_{α} . We shall simply say that ν on X is invariant if it satisfies this condition. Thus we can speak of the type of a measure class $C(\mu)$ with μ constructed as above without reference to G (just as long as G is built from some groups G_{α}). This also means that we are free to take any G_{α} we want in our proofs. For instance if X_{α} has n_{α} points ($n_{\alpha} = \infty$, possibly), then we can identify X_{α} and G_{α} with the integers mod n_{α} with G_{α} acting on itself by translation.

In general it will be convenient notationally to number the points of each X_{α}

by the integers from 1 to n_{α} . The choice of the numbering is, however, at our disposal. The data we are given then essentially are the measures μ_{α} , which in turn are determined by $p_{\alpha}^{i} = \mu_{\alpha}(\{i\})$ for the points i of X_{α} . Since $\sum_{i} p_{\alpha}^{i} = 1$ for every α , we can number the points of X_{α} so that p_{α}^{i} is decreasing as a function of i.

We shall now state the main results.

THEOREM 1. (1) $C(\mu)$ is type I if and only if $\sum_{\alpha} (1 - p_{\alpha}^{1})$ converges.

(2) $C(\mu)$ is type II_1 if and only if $n_{\alpha} < \infty$ for all α and $n_{\alpha} > 1$ for infinitely many α , and

(1)
$$\sum_{\alpha,i} (n_{\alpha})^{-1} |1 - (p_{\alpha}^{i} n_{\alpha})^{1/2}|^{2}$$

converges.

Bures [2] has obtained part (1) which is rather easy, and has obtained something equivalent to (2) under the additional hypothesis that the n_{α} are uniformly bounded. Our proof does not need this side condition. It is somewhat more subtle to distinguish the II_{α} from the III case, and we must impose a side condition similar but weaker than Bures' condition (but not 'essentially' weaker).

Let us establish the notation $|x|_c = \inf(|x|, c)$ for c > 0.

THEOREM 2. Suppose that $p_{\alpha}^1 \geq a > 0$ for some a and all α . Then $C(\mu)$ is type III if and only if

(2)
$$\sum_{\alpha,i} p_{\alpha}^{i} |(p_{\alpha}^{1}/p_{\alpha}^{i}) - 1|_{c}^{2}$$

diverges for some (and hence all) positive c.

Therefore, of course, we can distinguish the II_{∞} case at least under our side condition. Notice that $p_{\alpha}^1 \geq a > 0$ is automatically satisfied if the n_{α} are uniformly bounded by some K, for we can take a = 1/K. It is an easy matter to see that the convergence or divergence of the series in the theorem is independent of the value of c. We have stated our results in terms of absolutely convergent double series. These criteria can be rephrased in terms of conditionally convergent double series as in [2], but we feel they are more convenient in the present form.

PROOF OF THEOREM 1. We include the trivial proof of part 1, theorem 1 for completeness. Recall that $C(\mu)$ is type I if and only if μ gives positive measure to some point. If $x \in X$ with $x = (j_{\alpha})$, $1 \le j_{\alpha} \le n_{\alpha}$, then $\mu(\{x\}) = \prod \mu_{\alpha}(\{j_{\alpha}\})$, since μ is a product measure. Now it is clear that this product is as large as possible when $j_{\alpha} = 1$ for all α . Thus $C(\mu)$ is type I if and only if $\prod_{\alpha} p_{\alpha}^{1} > 0$, and this is true if and only if $\sum (1 - p_{\alpha}^{1})$ is convergent.

We pass now to part (2). Suppose that $C(\mu)$ is type II₁. Then by lemma 2.2, it is clear that n_{α} must be finite. Furthermore, if $n_{\alpha} = 1$ for almost all α , X is finite, and we are in the type I case. We must prove now that the series in question converges.

Let λ be a finite invariant measure on X. As we mentioned before, we may as well take X_{α} to the integers mod n_{α} and $G_{\alpha} = X_{\alpha}$ acting on itself. In that case X is a compact abelian group, and G may be viewed as a dense subgroup. Since λ is finite and invariant under translation by G, it is invariant under X. (To see this, look at the Fourier-Stieltjes transform of ν .) Therefore λ is the Haar

measure on X, and we may take $\lambda(X)=1$. Then if λ_{α} is the measure on X_{α} giving mass $(n_{\alpha})^{-1}$ to each point, it is clear that λ is the infinite product of the λ_{α} . Now since $\lambda \sim \mu = \prod \mu_{\alpha}$, we may apply a criterion of Kakutani [5] for equivalence of infinite product measures. His result says that $\lambda \sim \mu$ if and only if $\sum_{\alpha} |1 - (d\mu_{\alpha}/d\lambda_{\alpha})^{1/2}|^2$ converges, where $|\cdot|$ indicates the norm in $L_2(\lambda_{\alpha})$. But now this series can be seen by substituting the p_{α}^i , to be precisely the series in the statement of the theorem, and so we are done. (We are indebted to Kakutani for calling our attention to his paper.)

Conversely, let the conditions of the theorem be satisfied; then in the same notation, $\lambda \sim \mu$ so that $\lambda \in C(\mu)$, using Kakutani's result the other way. Thus as λ is invariant, $C(\mu)$ is either type I or II₁. The condition $n_{\alpha} > 1$ for infinitely many α assures us that we are in the II₁ case. This completes the proof.

3. Proof of first part of theorem 2

We come now to the proof of theorem 2. We first prove the sufficiency of our condition which is the more difficult part. If (a_n) and (b_n) are sequences of non-negative numbers, we say that $a_n \sim b_n$ if there exist positive constants k_1 and k_2 such that $k_1a_n \leq b_n \leq k_2a_n$. If $a_n \sim b_n$, then clearly $\sum a_n$ converges or diverges with $\sum b_n$.

Suppose now that $C(\mu)$ contains an invariant measure ν , and let $f = d\nu/d\mu$ so that f is a Borel function on X, unique up to μ null sets and > 0 almost everywhere. According to lemma 2.2, we have $\nu = \nu_{\alpha} \times \nu_{\alpha}^{*}$, and since $\mu = \mu_{\alpha} \times \mu_{\alpha}^{*}$, we see that $d\nu/d\mu = d\nu_{\alpha}/d\mu_{\alpha} \times d\nu_{\alpha}^{*}/d\mu_{\alpha}^{*}$. To be precise, if $x \in X$, $x = (x_{\alpha}x_{\alpha}^{*})$, then $f(x) = f_{\alpha}(x_{\alpha})f_{\alpha}^{*}(x_{\alpha}^{*})$ for almost all pairs $(x_{\alpha}, x_{\alpha}^{*})$ where $f_{\alpha} = d\nu_{\alpha}/d\mu_{\alpha}$ and $f_{\alpha}^{*} = d\nu_{\alpha}^{*}/d\mu_{\alpha}^{*}$. Since X_{α} is discrete, this holds for almost all x_{α}^{*} for each $x_{\alpha} \in X_{\alpha}$. Moreover, since ν_{α}^{*} is unique in this decomposition, it follows that ν_{α}^{*} is an invariant measure on X_{α}^{*} in the class $C(\mu_{\alpha}^{*})$. We can repeat this argument to find the following.

LEMMA 3.1. Let F be a finite subset of A, and let $X_F = \prod_{\alpha \in F} X_{\alpha}$ and $X_F^* = \prod_{\alpha \notin F} X_{\alpha}$. Then $X = X_F \times X_F^*$ where $x = (x_F, x_F^*)$ are the coordinates of a point x. Then $\nu = (\prod_{\alpha \in F} \nu_{\alpha}) \times \nu_F^*$ for some invariant measure ν_F^* on X_F^* and $f(x) = (\prod_{\alpha \in F} f_{\alpha}(x_{\alpha})) \cdot f_F^*(x_F^*)$ for almost all x_F^* where $f_F^* = d\nu_F^*/d\mu_F^*$ and $\mu_F^* = \prod_{\alpha \notin F} \mu_{\alpha}$.

Now let g^{λ} , g^{λ}_{α} , $g^{*\lambda}_{F}$ denote respectively the functions exp $(2\pi i\lambda \log (\cdot))$, where (\cdot) is f, f_{α} , f^{*}_{F} for real λ . These functions are of modulus one and

(3)
$$g^{\lambda}(x) = \prod_{\alpha \in F} g^{\lambda}_{\alpha}(x_{\alpha}) g^{*\lambda}_{F}(x_{F}^{*}).$$

We view each of these functions as an element of $L_2(X, \mu)$ by making them constant in the variables on which they do not depend. Since $\mu(X) = 1$, they are all unit vectors.

Now denote by P_F the orthogonal projection of $L_2(X)$ onto $L_2(X_F)$ where we view $L_2(X_F)$ as the subspace of $L_2(X)$ of all functions which depend only on the coordinates $x_{\alpha}, \bar{x}_{\alpha} \in F$, of the point $x \in X$. It is immediate that

$$(4) P_F(g^{\lambda}) = \left(\prod_{\alpha \in F} g_{\alpha}^{\lambda}\right) \cdot (g_F^{*\lambda}, 1)$$

where 1 denotes the function equal to one everywhere. Then it is well known that $P_F(g^{\lambda}) \to g^{\lambda}$ in the norm topology on $L_2(X)$. Therefore,

$$|P_F(g^{\lambda})| = |(g_F^{*\lambda}, 1)| \rightarrow 1,$$

and in particular $(g_F^{*\lambda}, 1) \neq 0$ for large F. It follows at once that $(g_{\alpha}^{\lambda}, 1) = 0$ can hold for only a finite number of α . Now let us choose constants $a_{\alpha}^{\lambda} = \exp(2\pi i \lambda b_{\alpha}^{\lambda})$ of modulus one such that $(a_{\alpha}^{\lambda} g_{\alpha}^{\lambda}, 1) \geq 0$. Of course b_{α}^{λ} is not determined at all for a finite set of α , and in general is only determined modulo $1/\lambda$.

Now we write

(5)
$$P_{F}(g^{\lambda}) = \left(\prod_{\alpha \in F} a_{\alpha}^{\lambda} g_{\alpha}^{\lambda}\right) \cdot \left(\prod_{\alpha \in F} a_{\alpha}^{-\lambda}\right) (g_{F}^{*\lambda}, 1),$$

and denote by γ_F the product of the last two terms. We observe that by our construction, for large F, the argument of γ_F independent of F. If we put $d = \bar{\gamma}_F/|\gamma_F|$ for large F, then

(6)
$$P_F(dg^{\lambda}) = \prod_{\alpha \in F} a_{\alpha}^{\lambda} g_{\alpha}^{\lambda} |\gamma_F|.$$

We have observed before that $|\gamma_F| \to 1$ as $F \to \infty$, and we immediately conclude that $\prod_{\alpha \in F} a_{\alpha}^{\lambda} g_{\alpha}^{\lambda}$ converges to dg^{λ} in $L_2(X)$.

On the other hand, one knows precisely when this happens.

LEMMA 3.2. Let $h_{\alpha} \in L_2(X_{\alpha})$ with $|h_{\alpha}| = 1$, and view them as elements of $L_2(X)$. Suppose that $(h_{\alpha}, 1) \geq 0$. Then $\prod_{\alpha \in F} h_{\alpha}$ converges in $L_2(X)$ if and only if $\sum |h_{\alpha} - 1|^2$ converges.

PROOF. This is well known (cf. [10], lemma 3.3.4, p. 24).

We apply this to the situation at hand to conclude that $\sum_{\alpha} |a_{\alpha}^{\lambda} g_{\alpha}^{\lambda} - 1|^2$ converges for every λ . Now we know what f_{α} , and hence g_{α}^{λ} , are; indeed $f_{\alpha}(i) = (d\nu_{\alpha}/d\mu_{\alpha})(i) = (p_{\alpha}^{j})^{-1}$ for $1 \leq j \leq n_{\alpha}$. We find then by putting in these values for g_{α}^{λ} that the series of nonnegative terms

(7)
$$\sum_{\alpha,i} p_{\alpha}^{i} |\exp(2\pi i \lambda (b_{\alpha}^{\lambda} - \log p_{\alpha}^{i})) - 1|^{2}$$

converges for all real λ . In particular, we may rearrange at will and conclude that the sum of the terms with i=1 converges. Since p^1_{α} is assumed bounded below, we conclude that the series

(8)
$$\sum_{\alpha} |\exp(2\pi i \lambda (b_{\alpha}^{\lambda} - \log p_{\alpha}^{1})) - 1|^{2}$$

converges for all λ . In particular the individual terms of the series approach zero. Now let ||x|| for a real number x denote the distance from x to the nearest integer. Then it is clear that for large α , the series (8) will dominate term by term the series

(9)
$$\sum_{\alpha} \|\lambda(b_{\alpha}^{\lambda} - \log p_{\alpha}^{1})\|^{2},$$

and hence the latter series converges. As we remarked before, the b_{α}^{λ} are deter-

mined only modulo $1/\lambda$, and now we make use of this freedom of choice to choose b_{α}^{λ} so that

(10)
$$\sum_{\alpha} |b_{\alpha}^{\lambda} - \log p_{\alpha}^{1}|^{2}$$

converges. This says that the b_{α}^{λ} are in some way almost independent of λ in the sense that they do not differ very much from a sequence which is independent of λ . Note that the hypothesis $p_{\alpha}^{1} \geq a > 0$ was used in a very crucial way to accomplish this.

Now if $s_{\alpha}^{\lambda} = \exp\left(-2\pi i\lambda(b_{\alpha}^{\lambda} - \log p_{\alpha}^{1})\right)$, then by (10), $\sum_{\alpha} |1 - s_{\alpha}^{\lambda}|^{2}$ converges. Lemma 3.3. Let the unit vectors h_{α} be as in lemma 3.2 with $(h_{\alpha}, 1) \geq 0$ and $\sum |h_{\alpha} - 1|^{2} < \infty$. If $\sum |s_{\alpha} - 1|^{2}$ converges with $|s_{\alpha}| = 1$, then $\sum |s_{\alpha}h_{\alpha} - 1|^{2}$ converges.

PROOF. This is quite obvious by the inequality $|x+y|^2 \le 2(|x|^2 + |y|^2)$ for vectors x and y.

We now apply lemma 3.3 to the vectors $a_{\alpha}^{\lambda}g_{\alpha}^{\lambda}$ and to s_{α}^{λ} . We then conclude the following, using the fact that $|\exp(ix) - 1|^2 = 2(1 - \cos x)$.

LEMMA 3.4. The series

(11)
$$\sum_{\alpha,i} p_{\alpha}^{i} (1 - \cos(2\pi\lambda \log(p_{\alpha}^{i}) - \log(p_{\alpha}^{i})))$$

converges for all real λ .

The following lemma now essentially completes the proof of the first part of theorem 2.

Recall that $|x|_c = \inf(|x|, c)$ for real x and c > 0.

Lemma 3.5. Let $\{a_n\}$ and $\{b_n\}$ be sequences of nonnegative numbers. Then the following are equivalent:

- (1) $\sum a_n(1-\cos(\lambda b_n))$ converges for all λ .
- (2) $\sum a_n(1-\cos(\lambda b_n))$ converges for λ in a set of positive Lebesgue measure.
- (3) $\sum a_n |b_n|_c^2$ converges for some c > 0.
- (4) $\sum a_n |b_n|_c^2$ converges for all c > 0.

Assuming this, let us complete the proof of the first part of theorem 2. According to lemmas 3.4 and 3.5, we see that $\sum_{\alpha,i} p_{\alpha}^{i} |\log p_{\alpha}^{1}/p_{\alpha}^{i}|_{c}^{2}$ converges for some, and hence all c > 0. Now observe that $|\log (p_{\alpha}^{1}/p_{\alpha}^{i})|_{c} \sim |(p_{\alpha}^{1}/p_{\alpha}^{i}) - 1|_{c}$, and hence we see that

(12)
$$\sum_{\alpha,i} p_{\alpha}^{i} |(p_{\alpha}^{1}/p_{\alpha}^{i}) - 1|_{c}^{2}$$

converges for some, and hence all c > 0. This proves the sufficiency of the condition in theorem 2.

PROOF OF LEMMA 3.5. Clearly $(1) \Rightarrow (2)$, and also $(3) \Rightarrow (4)$, for it is obvious that $a_n|b_n|_c^2 \sim a_n|b_n|_d^2$ for any two c, d > 0. Therefore, $\sum a_n|b_n|_c^2$ converges if and only if $\sum a_n|b_n|_d^2$ converges. Furthermore, $(4) \Rightarrow (1)$ is clear since, for fixed λ and c, a constant multiple of $a_n|b_n|_c^2$ dominates $a_n(1 - \cos(\lambda b_n))$ term by term. It remains then to prove that $(2) \Rightarrow (3)$.

Let S be a set of positive Lebesgue measure where the series in (2) converges,

and let $M(\lambda)$ be its sum for $\lambda \in S$. Clearly $M(\lambda)$ is a Borel function of λ , and consequently, by regularity of Lebesgue measure, we can find a compact set T of positive measure contained in S such that $M(\lambda) \leq K$ for some K and all λ in T. Then

(13)
$$\int_{T} M(\lambda) d\lambda \leq K|T| < \infty,$$

where |T| is the Lebesgue measure of T. Since we are dealing with a series of positive terms, we can invert the order of summation and integration to find that

(14)
$$\sum a_n \int_T (1 - \cos(\lambda b_n)) d\lambda$$

converges. Now we claim that the function $f(t) = \int_T (1 - \cos(\lambda t)) d\lambda$ has a lower bound m(c) on the set $|t| \ge c$. In fact f(t) is continuous in t, and positive except for t = 0. It suffices to show then that it is bounded from zero at infinity. But this is clear for as $|t| \to \infty$, $f(t) \to |T|$ by the Riemann-Lebesgue lemma.

Therefore, if $L(c) = \{n: |b_n| \ge c\}$, we see that $\sum_{n \in L(c)} a_n \le (m(c))^{-1}K|T| < \infty$. On the other hand, by assumption $\sum a_n(1 - \cos \lambda b_n)$ converges for some $\lambda > 0$. Then choose c so small that $1 - \cos x \ge x^2/4$ if $|x| < c\lambda$. Then

(15)
$$\sum a_n |b_n|_c^2 = \sum_{n \in L(c)} a_n |b_n|_c^2 + \sum_{n \in L(c)} a_n,$$

and the first series is dominated term by term by $4 \sum a_n (1 - \cos \lambda b_n)$, $(n \notin L_c)$, and the second series is convergent as we have shown above. This completes the proof.

4. Restricted products of measures

We now have proved in theorem 2 that if $C(\mu)$ is not of type III, then the series of the theorem converges. To show the converse, we must be able to construct invariant measures λ on $X = \prod X_{\alpha}$, provided that the infinite series of the theorem converges. This section is then devoted to a short digression on a method of construction of measures on product spaces which will have applications beyond the immediate scope of this paper.

Suppose that S_{α} is a standard Borel space for each $\alpha \in A$, a countable index set, and let K_{α} be a Borel subset of S_{α} . We let S be the full Cartesian product of the S_{α} with the product Borel structure, and we let $E = \{(s_{\alpha}): s_{\alpha} \in K_{\alpha} \text{ almost everywhere}\}$. We call E the restricted product of the spaces S_{α} relative to the subsets K_{α} .

LEMMA 4.1. The restricted product E is a Borel subset of S, and hence a standard Borel space; moreover, $E = S_{\alpha} \times E_{\alpha}^*$ where E_{α}^* is the restricted product of the S_{β} , $\beta \neq \alpha$.

PROOF. If F is a finite subset of A, let $S^F = \prod_{\alpha \in F} S_\alpha \times \prod_{\alpha \notin F} K_\alpha$. Then S^F is clearly a Borel subset of F. Moreover, $E = \bigcup_F S^F$ so that E is also a Borel subset of S. The second statement is obvious.

We note that if σ is a subset of E, then σ is a Borel subset of E if and only if

 $\sigma \cap S^F$ is a Borel subset of S^F for every finite set F. Thus the Borel structure on E is the 'inductive limit' of the Borel structures on S^F . This method could also be used to define the Borel structure on E independently of its structure as a subset of S. Suppose further that λ_{α} is a family of measures on S_{α} such that $\lambda_{\alpha}(K_{\alpha}) = 1$. We want to define the restricted product measure $\lambda = \prod \lambda_{\alpha}$ as a measure on E.

Indeed, let F be a finite set of indices, and let λ^F be the product of the measures λ_{α} ($\alpha \in F$) and the measures λ_{α} restricted to K_{α} , $\alpha \notin F$ on $S^F = \prod_{\alpha \in F} S_{\alpha} \times \prod_{\alpha \notin F} K_{\alpha}$. This is well defined since $\lambda_{\alpha}(K_{\alpha}) = 1$, and is a (σ -finite) measure on S^F . Then if σ is a Borel subset of E, define $\lambda(\sigma) = \lim \lambda^F(\sigma \cap S^F)$. The limit exists as the sequence on the right increases as F increases.

LEMMA 4.2. The restricted product measure λ is a $(\sigma$ -finite) measure on E and $\lambda = \lambda_{\alpha} \times \lambda_{\alpha}^{*}$ (product measure), where $E = S_{\alpha} \times E_{\alpha}^{*}$ and λ_{α}^{*} is the restricted product of the λ_{β} , $\beta \neq \alpha$ on E_{α}^{*} .

Proof. The assertions are completely obvious.

Since E is a subset of S, we may always regard any measure λ on E as a measure on S by $\lambda(\sigma) = \lambda(\sigma \cap E)$ for Borel subsets σ of S. With this convention the following is clear.

COROLLARY. Let X_{α} (= S_{α}) be discrete, and let K_{α} be a finite subset of X_{α} consisting of m_{α} elements. Let λ_{α} be the measure on X_{α} giving mass $(m_{\alpha})^{-1}$ to each point, so that $\lambda_{\alpha}(K_{\alpha}) = 1$. Then the restricted product measure λ on $X = \prod X_{\alpha}$ is invariant (cf. lemma 2.2) and ergodic.

Proof. This follows from the definition of invariance and the second statement of the previous lemma. The ergodicity follows by an easy modification of lemma 2.1.

5. Completion of the proof of theorem 2

We now return to the second half of the proof of theorem 2. Thus we assume that the series

(16)
$$\sum_{\alpha,i} p_{\alpha}^{i} |(p_{\alpha}^{1}/p_{\alpha}^{i}) - 1|_{c}^{2}$$

converges for all c > 0, and that $p_{\alpha}^1 \ge a > 0$ for some a. It is convenient to extend the definition of p_{α}^i for $i > n_{\alpha}$ (= number of points in X_{α}) by making it zero for $i > n_{\alpha}$.

We first notice a trivial remark. As before, if F is a finite set of indices, let $X_F = \prod_{\alpha \in F} X_{\alpha}$ and $X_F^* = \prod_{\alpha \notin F} X_{\alpha}$, and similarly for μ_F and μ_F^* .

Lemma 5.1. The class $C(\mu)$ on X is type III if and only if μ_F^* on X_F^* is type III. Proof. If $\nu \in C(\mu)$ is invariant, it follows by repeated application of lemma 2.2 that $\nu = \nu_F \times \nu_F^*$, where ν_F gives each point of X_F mass one and $\nu_F^* \in C(\mu_F^*)$ is invariant. Conversely, if $\nu_F^* \in C(\mu_F^*)$ is invariant, then ν defined by the above formula is in $C(\mu)$, and is invariant. This completes the proof.

The meaning of this is that we are free to delete any finite set of indices from A

as far as the question of type III is concerned. Since the convergence or divergence of the series in theorem 2 clearly does not depend on any finite set F of indices from A, we may ignore finite sets of indices when proving theorem 2.

Now if b > 0 is fixed and k is a positive integer, define a subset L(b, k) of A by $L(b, k) = \{\alpha : p_{\alpha}^{k} \leq b\}$. Since $p_{\alpha}^{k+1} \leq p_{\alpha}^{k}$, it is clear that $L(b, k) \subset L(b, k+1)$. If we take b = a/2 where $p_{\alpha}^{1} \geq a$, then

(17)
$$\sum_{\alpha}' p_{\alpha}^{i} |(p_{\alpha}^{1}/p_{\alpha}^{i}) - 1|_{1}^{2} = \sum_{\alpha}' p_{\alpha}^{i} < \infty,$$

where the primed summation is over all $\alpha \in L(b, k)$ and all $i \geq k$. On the other hand, if $\alpha \notin L(b, k)$, and $i \leq k$, then $p_{\alpha}^{i} \geq b > 0$, and so

(18)
$$\sum |(p_{\alpha}^{1}/p_{\alpha}^{i})-1|_{c}^{2}, \qquad (\alpha \notin L(b,k), i \leq k)$$

converges for any c > 0. This of course immediately implies that

(19)
$$\sum |p_{\alpha}^{1}/p_{\alpha}^{i}-1|^{2}, \qquad (\alpha \notin L(b,k), i \leq k)$$

converges since $p_{\alpha}^{i}/p_{\alpha}^{i}$ is bounded on this range of indices. Now suppose 1/n < a. Then the convergence of (19) implies that except for a finite number of $\alpha \notin L(b, k)$, we have $p_{\alpha}^{i} \geq 1/n$ for all $i \leq k$. Therefore, if A - L(b, k) is infinite, there exists an α such that $p_{\alpha}^{i} \geq 1/n$ for all $i \leq k$. It follows then that $1 \geq \sum_{i \leq k} p_{\alpha}^{i} \geq k/n$. In other words, if k > n, then A - L(b, k) must be a finite set. By our previous remarks, we are free to assume that A = L(b, n + 1).

We now define a partition of A into n disjoint sets A(k), $1 \le k \le n$ by A(k) = L(b, k + 1) - L(b, k) so that $\alpha \in A(k)$ if and only if $p_{\alpha}^{k} > b$, and $p_{\alpha}^{k+1} \le b$. It is clear that if $\alpha \in A(k)$, then X_{α} has at least k points. We deduce then from (17) and (19) the convergence of

(20)
$$\sum p_{\alpha}^{i}, \qquad (\alpha \in A(k), i > k)$$

and of

(21)
$$\sum |(p_{\alpha}^{1}/p_{\alpha}^{i}) - 1|^{2}, \qquad (\alpha \in A(k), i \leq k).$$

Now let K_{α} be the subset of X_{α} consisting of the first k points when $\alpha \in S(k)$. We let ν_{α} be the measure on X_{α} giving mass 1/k to each point when $\alpha \in S(k)$, and we let ν be the restricted product of the ν_{α} using the sets K_{α} . Then, according to the corollary to lemma 4.2, ν is an invariant measure on X. The following completes the proof of theorem 2.

LEMMA 5.2. Let μ and ν be as above, then $\nu \in C(\mu)$.

PROOF. Since ν and μ are both ergodic under some group G acting on X, they are either equivalent or mutually singular. Thus it suffices to show that μ and ν are not mutually singular. Let K be the full Cartesian product of the sets K_{α} . Then $\nu(K) = 1$ by construction. On the other hand, if $a_{\alpha} = \sum_{i>k(\alpha)} p_{\alpha}^{i}$ where $k(\alpha) = k$ if $\alpha \in S(k)$, then by (20), $\sum a_{\alpha}$ converges. But now $\mu(K) = \prod_{\alpha} \mu_{\alpha}(K_{\alpha}) = \prod_{\alpha} (1 - a_{\alpha})$, and this infinite product is therefore positive. Now let λ_{α} be the measure on K_{α} which is $(1 - a_{\alpha})^{-1}$ times the restriction of μ_{α} to K_{α} . Then $\lambda_{\alpha}(K_{\alpha}) = 1$ and the infinite product of the λ_{α} , say λ , is a multiple of μ restricted to K. Thus to complete the proof, it is enough to show that

 $\lambda \sim \nu$ as measures on K. Observe that λ_{α} gives mass $(1 - a_{\alpha})^{-1}p_{\alpha}^{i} = q_{\alpha}^{i}$ to the i-th point of K_{α} , $i \leq k(\alpha)$. Then by (21),

$$\sum |(q_{\alpha}^1/q_{\alpha}^i) - 1|^2$$

converges where the sum extends over all α and all $i \leq k(\alpha)$.

Now since q_{α}^{t} has a lower bound, b, over the range of summation, we see from the above that

(23)
$$\sum (q_{\alpha}^{1} - q_{\alpha}^{i})^{2}, \qquad (\alpha, i \leq k(\alpha))$$

converges. Now using the fact that $\sum_{j=1}^{k(\alpha)} q_{\alpha}^{j} = 1$, and the Schwarz inequality, we deduce immediately that

(24)
$$k(\alpha)(q_{\alpha}^{1} - k(\alpha)^{-1})^{2} \leq \sum_{i=1}^{k(\alpha)} (q_{\alpha}^{1} - q_{\alpha}^{i})^{2}.$$

Now summing on α , and using (23) and the inequality $(s+t)^2 \leq 2(s^2+t^2)$, we see that

(25)
$$\sum \left((q_{\alpha}^{1}/q_{\alpha}^{i}) - (k(\alpha)q_{\alpha}^{i})^{-1} \right)^{2}, \qquad (\alpha, i \leq k(\alpha))$$

is dominated term by term by a convergent series, and so also converges. Now combining (22) and (25) and the same numerical inequality, we see that

(26)
$$\sum ((q_{\alpha}^{i}k(\alpha))^{-1} - 1)^{2}, \qquad (\alpha, i \leq k(\alpha))$$

is dominated term by term by twice the sum of (22) and (25), and thus converges. Finally we observe that $(q^i_{\sigma}k(\alpha))$ is bounded above and below so that

$$((q_{\alpha}^{i}k(\alpha))^{-1}-1)^{2} \sim ((q_{\alpha}^{i}k(\alpha))^{1/2}-1)^{2};$$

thus,

(28)
$$\sum_{\alpha,j \leq k(\alpha)} ((q_{\alpha}^{i}k(\alpha))^{1/2} - 1)^{2}$$

is convergent. Now by Kakutani's criterion [5] for equivalence of infinite product measures used in section 2, we find that $\lambda \sim \nu$, and this completes the proof of the lemma and hence of theorem 2.

6. Complements

We would like to add a few words about the meaning of the series $\sum p_{\alpha}^{i}(p_{\alpha}^{1}/p_{\alpha}^{i}-1)^{2}$ of theorem 2. For this to converge, either p_{α}^{i} must be small, or the second term must be small. In other words, the convergence of this series means that asymptotically the measure μ_{α} gives approximately equal mass to some set of points in X_{α} (those i such that the second factor is small) and very small mass to the rest. Note that both factors cannot both be simultaneously small; this is the meaning of the hypothesis $p_{\alpha}^{1} \geq a > 0$. We can interpret the condition for $C(\mu)$ to be type II₁ as simply that, all the points of X_{α} get approximately the same mass, and the type I condition can be interpreted to say that the first point in X_{α} carries nearly all the mass with the other points having very small mass.

We do not know what happens when the condition $p_{\alpha}^{1} \geq a > 0$ is dropped. However, one can show by the same arguments as used in the second part of the proof of theorem 2, that if there exist integers $k(\alpha)$ such that

(29)
$$\sum_{\alpha,i>k(\alpha)} p_{\alpha}^{i}$$

and

(30)
$$\sum_{\alpha,i \leq k(\alpha)} |(k(\alpha))^{-1/2} - (p_{\alpha}^i)^{1/2}|^2$$

converge, then $C(\mu)$ is not type III.

Let us consider for a moment an example of theorem 2 which was our original motivation. Let A be the set of rational primes, and let X_q be infinite for each prime q with μ_q defined by $p_q^i = q^{1-i}(1-q^{-1})$ for $i=1,2,\cdots$. Then since $p_q^1 \geq \frac{1}{2}$ for all q, theorem 2 applies. We observe that terms of the series in theorem 2 for i=2 are asymptotic to q^{-1} , and since $\sum q^{-1}$ diverges, we are in the type III case. It is amusing to note that we do not need lemma 3.5 to do this. Indeed if we look at the series in lemma 3.4 instead, and take the terms with i=2, we see that it is termwise asymptotic to

(31)
$$\sum q^{-1}(1-\cos(2\pi\lambda\log q)),$$

and if we were in the type I or II case, this would converge for all real λ . But it is well known that

(32)
$$\sum q^{-1} \cos (2\pi\lambda \log q)$$

is (conditionally) convergent for real $\lambda \neq 0$ (see [12], p. 56); this is equivalent to the fact that the Dirichlet series for $\log |\zeta(s)|$ converges for Re (s) = 1, $s \neq 1$, and hence that the zeta-function, $\zeta(s)$, has no zeroes on this line). Now by addition, we see that $\sum q^{-1}$ would converge which is of course nonsense.

We note that this example is not covered by Bures' criterion ([2], p. 170) for the type III case. This example will be used elsewhere to show that the adele group of the 'ax + b' group [14] has as regular representation a type III factor representation.

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