# Approximately transitive flows and ITPFI factors 

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This paper is dedicated to Richard V. Kadison on the occasion of his sixtieth birthday.


#### Abstract

We define a new property of a Borel group action on a Lebesgue measure space, which we call approximate transitivity. Our main results are (i) a type $\mathrm{III}_{0}$ hyperfinite factor is ITPFI if and only if its flow of weights is approximately transitive, and (ii) for ergodic transformations preserving a finite measure, approximate transitivity implies zero entropy.


## 0 . Introduction

von Neumann algebras are the non-commutative analogue of measure theory spaces. The product measures of measures on finite sets give rise to a class of factors called ITPFI factors (see terminology). However, in the classification problem the most natural class turned out to be the approximately type I factors [4] (those factors which are well-approximated by finite-dimensional ones - see terminology). It is trivial that ITPFI implies approximately type I, but the converse is false and non-trivial [12], [2], [6]. (Since all non-atomic standard Borel measures are Borel isomorphic, the corresponding problem does not arise in measure theory.) The ITPFI factors are certainly the most natural subclass of the approximately type I factors. Their exact position among the approximately type I factors (up to isomorphism of factors) has remained an interesting mystery for some time. In particular, there is still no direct spatial construction of a non-ITPFI approximately type I factor. The crucial existential step is always carried out in the flow of weights (an ergodic flow which is naturally defined as an invariant of the factor). The problem only arises for factors of type $\mathrm{III}_{0}$ or $\mathrm{III}_{1}$. The known examples of non-ITPFI approximately type I factors are all of type $\mathrm{III}_{0}$. In this paper we completely characterize the ITPFI factors among the type $\mathrm{III}_{0}$ approximately type I factors by a new ergodic property of their complete invariant, the flow of weights, which we
call approximate transitivity. Of course this transfers the original problem to understanding approximate transitivity for ergodic flows. Our second major result is that for finite measure preserving flows, approximate transitivity implies zero entropy.

Approximate transitivity is a new and apparently interesting notion in ergodic theory. Our proof that ITPFI implies approximate transitivity is rather straightforward (lemma 8.1) and is, in fact, how this property was discovered. Our proof of the converse is of particular interest because it is obtained by attempting directly to 'invert the flow of weights arrow'. The argument is quite similar to the Murray-von Neumann proof of the uniqueness of the hyperfinite $\mathrm{II}_{1}$ factor [15]. They embed a finite-dimensional algebra in a finite type I factor in a very precise manner relative to the trace. For type III factors no trace exists and, instead of only comparing minimal projections inside the finite-dimensional algebra, we compare such minimal projections together with the restriction of a state to these projections. This comparison directly yields measure theoretical objects on the flow of weights, and (except in the $\mathrm{III}_{1}$ case) is non-trivial even though the comparison of projections is trivial. Our paper is intended to illustrate this technique. In fact an alternate proof of Krieger's theorem, not using the cohomological technique of Krieger, can be based on the original Murray-von Neumann proof of the uniqueness of the hyperfinite $\mathrm{II}_{1}$ together with the above refined comparison of projections. (Krieger's theorem [13] states, in part, that the flow of weights considered as a mapping from type III $_{0}$ Krieger factors with algebraic isomorphism as the equivalence relation, to strictly ergodic flows with conjugacy as the equivalence relation, is one-to-one and onto between equivalence classes.)

It is easy to translate our proof to the purely ergodic setting of non-singular transformations. Our result would then follow from Krieger's theorem (our proof does not use Krieger's theorem). However, as mentioned above, part of our goal was to exhibit the flow of weights as a useful technique. Indeed we present more of the comparison theory of finite weights than is needed for our proof.
§ 1 contains some terminology. In § 2 we define approximate transitivity (hereafter referred to as AT) for Borel group actions and give some elementary properties. In § 3 we prove that for finite measure preserving transformations, AT implies zero entropy. In § 4 we give three different constructions of the flow of weights which will be used later. $\S 5$ contains a comparison theory for finite periodic weights, and § 6 gives the comparison theory for finite (not necessarily periodic) weights. In § 7 we introduce a 'product property' which is equivalent to being ITPFI. In § 8 we prove the equivalence of the ITPFI and AT properties.

## 1. Terminology

A von Neumann algebra $M$ is said to be approximately type $I$ if it is of the form

$$
M=\left(\bigcup_{n=1}^{\infty} M_{n}\right)^{\prime \prime},
$$

where $M_{n} \subset M_{n+1}$ for each $n$, and each $M_{n}$ is a finite-dimensional matrix algebra (the names approximately finite, hyperfinite, approximately finite dimensional, and
matricial have all been used in the literature for this concept). A factor $M$ is said to be ITPFI if it is of the form $M=\otimes_{n=1}^{\infty}\left(M_{n}, \phi_{n}\right)$ where each $M_{n}$ is a finite type I factor ([1],[21]). A factor $M$ is called a Krieger factor if it can be obtained from an ergodic action of $\mathbb{Z}$ by the Murray-von Neumann group measure space construction. (It is straightforward that ITPFI $\Rightarrow$ Krieger $\Rightarrow$ approximately type I.)

For a detailed explanation of the following standard terminology see, for example [20]. The term weight always means a normal semi-finite weight. If $\phi$ is a weight on the von Neumann algebra $M$ then $\sigma_{i}^{\phi}$ is the modular automorphism group. The invariant $T(M)$ is the set of all $t$ such that $\sigma_{t}^{\phi}$ is inner. $M_{*}$ denotes the predual of $M$, and $M_{*}^{+}$is then the set of all finite weights. If $\psi \in M_{*}^{+}$then $s(\psi)$ denotes the support of $\psi$. The flow of weights of $M$ is an ergodic action of $\mathbb{R}_{+}^{*}$ on some measure space $\left(X_{M}, \mu_{M}\right)$. The construction of [5] gives not that measure space, but the measure algebra whose elements are equivalence classes [ $\Phi$ ] of integrable weights $\Phi$ of infinite multiplicity. The flow is then defined by $\mathscr{F}_{t}^{M}[\Phi]=[t \Phi]$. It is sometimes convenient to consider the flow as an action of $\mathbb{R}$, in which case it is written $F_{t}^{M}=\mathscr{F}_{e}^{M}$ (if $M$ is understood, it is usually omitted).

If $f$ is a function on a measure space $(X, \mu)$ then $\|f\|$ denotes the $L^{1}$-norm of $f$. If $\mu, \nu$ are finite measures on $X$ then $\|\mu-\nu\|$ is the $L^{1}$-norm defined by $\| d \mu / d \sigma-$ $d \nu / d \sigma \|$ where $\mu, \nu<\sigma$. If $x$ is a finite weight or operator then $\|x\|$ denotes the usual norm.

## 2. AT actions-elementary properties

We define approximate transitivity of a Borel group action on a Lebesgue measure space, and establish some elementary properties.
Definition 2.1. Let $G$ be a Borel group, $(X, \nu)$ a Lebesgue measure space, and $\alpha: G \rightarrow \operatorname{Aut}(X, \nu)$ a Borel homomorphism. We say that the action is approximately transitive (AT) if given $n<\infty$, finite measures $\mu_{1}, \ldots, \mu_{n}<\nu$, and $\varepsilon>0$, there exists a finite measure $\mu<\nu, g_{1}, \ldots, g_{m} \in G$ for some $m<\infty$, and $\lambda_{j k} \geq 0, k=1, \ldots, m$ such that

$$
\begin{equation*}
\left\|\mu_{j}-\sum_{k=1}^{m} \lambda_{j k} \alpha_{g_{k}} \mu\right\| \leq \varepsilon, \quad j=1, \ldots, n . \tag{2.1}
\end{equation*}
$$

If $G=\mathbb{Z}$ and $\alpha$ is AT, then we say that $T=\alpha(1)$ is AT.
Remark 2.2. There are a number of elementary variations on this definition.
(i) The index $k$ need not be restricted to a finite set. Typically we will take $k \in \mathbb{Z}$ and consider $\lambda_{j k}$ as a function $\lambda_{j} \in \ell_{+}^{1}(\mathbb{Z})$.
(ii) One can demand that $\|\mu\|=1$ and $\left\|\lambda_{j}\right\|=\sum_{k} \lambda_{j k}=\left\|\mu_{j}\right\|$.
(iii) By taking $\|\mu\|$ sufficiently small, one can take the $\lambda_{j k}$ to be integers.
(iv) It is sufficient to ask that eq. (2.1) hold for $n=2$. (If $\mu^{\prime}$ approximates $\mu_{1}, \ldots, \mu_{n-1}$ in the sense of eq. (2.1), choose $\mu$ to approximate $\mu^{\prime}$ and $\mu_{n}$ ).
(v) For continuous actions of a locally compact group, the $\lambda_{j k}$ can be replaced by functions $\lambda_{j} \in L_{+}^{1}(G, d g)$ such that

$$
\begin{equation*}
\left\|\mu_{j}-\int_{G} d g \lambda_{j}(g) \alpha_{g} \mu\right\| \leq \varepsilon, \quad j=1, \ldots, n . \tag{2.2}
\end{equation*}
$$

To show that eq. (2.2) implies eq. (2.1), approximate the $\lambda_{j}(g)$ by simple functions. To prove the converse, write the sum as an integral over delta functions $\delta\left(g g_{k}^{-1}\right) \lambda_{j k}$ and then approximate by functions in $L_{+}^{1}(G, d g)$.
(vi) The equation $h=d \sigma / d \nu$ gives a one-to-one correspondence between functions $h \in L_{+}^{1}(X, \nu)$ and finite measures $\sigma<\nu$. We have

$$
d\left(\alpha_{g} \sigma\right) / d \nu=\rho_{g} \alpha_{g}(d \sigma / d \nu)
$$

where $\rho_{g}=d\left(\alpha_{g} \nu\right) / d \nu$. Then

$$
\begin{equation*}
\left(\beta_{g} f\right)(x)=f\left(g^{-1} x\right) \rho_{\mathbf{g}}(x) \tag{2,3}
\end{equation*}
$$

defines a homomorphism $\beta$ from $G$ into the invertible isometries on $L^{1}(X, \nu)$. It is therefore equivalent to ask that for any $f_{1}, \ldots, f_{n} \in L_{+}^{1}(X, \nu)$ and $\varepsilon>0$, there exist $f \in L_{+}^{1}(X, \nu), g_{1}, \ldots, g_{m} \in G$ and $\lambda_{j k} \geq 0$ such that

$$
\begin{equation*}
\left\|f_{j}-\sum_{k=1}^{m} \lambda_{j k} \beta_{g_{k}} f\right\| \leq \varepsilon, \quad j=1, \ldots, n . \tag{2.4}
\end{equation*}
$$

If $\alpha_{g} \nu=\nu$ for all $g \in G$, then $\beta=\alpha$.
Lemma 2.3. An AT action is ergodic.
Proof. Let $B \subset X, \nu(B)>0, \nu(X \backslash B)>0$, and $\alpha_{g} B=B$ for all $g \in G$. Choose $B_{1} \subset B$ and $B_{2} \subset X \backslash B$ such that $0<\nu\left(B_{j}\right)<\infty, j=1$, 2. Let $\mu_{j}=\left.\nu\right|_{B_{j}}$. Then eq. (2.1) for $j=2$ implies that $\mu(B)<\varepsilon$. If $\varepsilon<\frac{1}{2} \nu\left(B_{1}\right)$ this contradicts eq. (2.1) for $j=1$.

Remark 2.4. Let ( $X, \mathscr{B}, \nu, G, \alpha$ ) be a Borel group action, $\nu(X)<\infty$. Let $\mathscr{B}_{0}$ be a sub- $\sigma$-algebra of the $\sigma$-algebra $\mathscr{B}$ of Borel subsets of $X$, such that $\alpha_{g} \mathscr{B}_{0}=\mathscr{B}_{0}$ for all $g \in G$. Then the restriction $\left(X, \mathscr{B}_{0}, \nu, G, \alpha\right)$ is called a factor action of the given action. If $\sigma$ is any finite measure on $(X, \mathscr{B})$ we have $\|\sigma\|_{\left(X, \mathscr{B}_{0}\right)} \leq\|\sigma\|_{(X, \mathscr{B})}$. Hence any factor action of an AT action is again AT.

The base and ceiling function construction of a flow is particularly useful when the ceiling function is constant. In this situation one naturally expects that the flow will have a certain property if and only if the base transformation has the corresponding property.

Lemma 2.5. Let ( $X, \nu, F_{s}$ ) be a flow built over the base transformation ( $B, \nu_{B}, T$ ) with a ceiling function with constant height $H$. Then $F_{s}$ is AT if and only if $T$ is AT. Proof. We can write $X=B \times I$ where the interval $I=[0, H)$ carries Lebesgue measure. We have

$$
\begin{equation*}
F_{s}(b, t)=\left(T^{n} b, u\right), \tag{2.5}
\end{equation*}
$$

where $s+t=n H+u, n \in \mathbb{Z}, 0 \leq u<H$.
Assume that $T$ is AT. Let $\mu_{1}, \ldots, \mu_{n}<\nu, \mu_{j}(X)<\infty$, and $\varepsilon>0$. Since rectangles generate the measure algebra, it follows by a straighforward but tedious argument that one can approximate the $\mu_{j}$ by a sum of product measures. More precisely, there exists an integer $L$ and measures $\mu_{j k}, k=0, \ldots, L-1$ on $B$ such that

$$
\begin{equation*}
\left\|\mu_{j}-\sum_{k=0}^{L-1} \mu_{j k} \times m_{k}\right\|<\varepsilon, \quad j=1, \ldots, n \tag{2.6}
\end{equation*}
$$

where $m_{k}$ is the restriction of Lebesgue measure to the interval $L_{k}=$ $[k H / L,(k+1) H / L)$. Since $T$ is AT there exist $\sigma_{B}<\nu_{B}$ and $\lambda_{j k} \in \ell_{+}^{1}(\mathbb{Z})$ such that

$$
\begin{equation*}
\left\|\mu_{j k}-\sum_{q \in \mathbb{Z}} \lambda_{j k}(q) T^{q} \sigma_{B}\right\|<\varepsilon L^{-1} . \tag{2.7}
\end{equation*}
$$

Define

$$
\begin{equation*}
\sigma=\sigma_{B} \times m_{0} \tag{2.8}
\end{equation*}
$$

and

$$
\begin{equation*}
\lambda_{j}(q L+k)=\lambda_{j k}(q), \quad q \in \mathbb{Z}, k=0, \ldots, L-1 . \tag{2.9}
\end{equation*}
$$

Then

$$
\begin{align*}
\left\|\mu_{j}-\sum_{r \in \mathbb{Z}} \lambda_{j}(r) F_{r H / L} \sigma\right\| & <\varepsilon+\sum_{k=0}^{L-1}\left\|\left.\left(\mu_{j k} \times \varepsilon_{k}-\sum_{r \in \mathbb{Z}} \lambda_{j}(r) F_{r H / L} \sigma\right)\right|_{B \times I_{k}}\right\| \\
& =\varepsilon+\sum_{k=0}^{L-1}\left\|\mu_{j k}-\sum_{q \in \mathbb{Z}} \lambda_{j k}(q) T^{q} \sigma_{B}\right\|<2 \varepsilon . \tag{2.10}
\end{align*}
$$

Thus $F_{s}$ is AT.
Now assume that $F_{s}$ is AT. Let $\mu_{1}, \ldots, \mu_{n}<\nu_{B}, \mu_{j}(B)=1$, and $0<\varepsilon<\frac{1}{2}$. Let $J_{k}=[k H / 6,(k+1) H / 6), k=0,1, \ldots, 5$, and let $m^{k}$ denote the restriction of Lebesgue measure to $J_{k}$. Let

$$
\begin{equation*}
\tilde{\mu}_{j}=\mu_{j} \times m^{3}, \quad j=1, \ldots, n . \tag{2.11}
\end{equation*}
$$

Then there exist $\mu<\nu,\|\mu\|=1$ and $\lambda_{j} \in \ell_{+}^{\prime}(\mathbb{Z}),\left\|\lambda_{j}\right\|=1$, and $s_{k} \in \mathbb{R}, k \in \mathbb{Z}$ such that

$$
\begin{equation*}
\left\|\tilde{\mu}_{j}-\sum_{k \in \mathbb{Z}} \lambda_{j}(k) F_{s_{k}} \mu\right\|<\varepsilon / 6, \quad j=1, \ldots, n . \tag{2.12}
\end{equation*}
$$

In particular we have

$$
\begin{equation*}
\left\|\left.\left(\sum_{k \in \mathbb{Z}} \lambda_{j}(k) F_{s_{k}} \mu\right)\right|_{B \times J_{3}^{J}}\right\|<\varepsilon / 6 . \tag{2.13}
\end{equation*}
$$

In order to produce the desired measure $\sigma$ on $B$ and $\Lambda_{j} \in \ell_{+}^{1}(\mathbb{Z})$ it is necessary to restrict the supports of $\lambda_{j}$ and $\mu$ somewhat. Since $\|\mu\|=1$ there is some $0 \leq K \leq 5$ such that

$$
\begin{equation*}
\left\|\left.\mu\right|_{B \times J_{K}}\right\| \geq \frac{1}{6} . \tag{2.14}
\end{equation*}
$$

By shifting the $s_{k}$ and shifting $\mu$ under $F_{s}$ (if necessary) we can assume that $K=4$. Let

$$
\begin{equation*}
Y=\left\{k \in \mathbb{Z}: 0 \leq s_{k}<\frac{1}{3} H(\bmod H)\right\} . \tag{2.15}
\end{equation*}
$$

Then $k \notin Y$ implies that $\left(F_{s_{k}}\left(\left.\mu\right|_{B \times J_{4}}\right)\right)\left(B \times J_{3}\right)=0$. It now follows from eqs. (2.13) and (2.14) that

$$
\begin{equation*}
\sum_{k \in Y} \lambda_{j}(k)<\varepsilon, \tag{2.16}
\end{equation*}
$$

and thus

$$
\begin{equation*}
\sum_{k \in Y} \lambda_{j}(k)>1-\varepsilon, \quad j=1, \ldots, n . \tag{2.17}
\end{equation*}
$$

Write $\mu=\mu^{\prime}+\mu^{\prime \prime}$ where

$$
\begin{equation*}
\mu^{\prime \prime}=\left.\sum_{k=0}^{2} \mu\right|_{B \times J_{k}} \tag{2.18}
\end{equation*}
$$

If $k \in Y$ then $\left(F_{s_{k}} \mu^{\prime \prime}\right)\left(B \times J_{3}\right)=0$, and eqs. (2.13) and (2.17) imply that

$$
\begin{equation*}
\left\|\mu^{\prime \prime}\right\|<(1-\varepsilon)^{-1} \varepsilon / 6<\frac{1}{3} \varepsilon \tag{2.19}
\end{equation*}
$$

Eqs. (2.17) and (2.19) imply that

$$
\begin{equation*}
\left\|\sum_{k \in \mathbb{Z}} \lambda_{j}(k) F_{s_{k}} \mu-\sum_{k \in Y} \lambda_{j}(k) F_{s_{k}} \mu^{\prime}\right\|<\frac{4}{3} \varepsilon \tag{2.20}
\end{equation*}
$$

Let $P_{*} \sigma$ denote the canonical projection of the finite measure $\sigma$ on $\dot{X}=B \times I$ onto $B$. We have

$$
\begin{equation*}
P_{*} \tilde{\mu}_{j}=(H / 6) \mu_{j} \tag{2.21}
\end{equation*}
$$

and

$$
\begin{equation*}
P_{*} F_{s_{k}} \mu^{\prime}=T^{p}\left(P_{*} \mu^{\prime}\right) \tag{2.22}
\end{equation*}
$$

if $s_{k} \in Y_{p}$ where

$$
\begin{equation*}
Y_{p}=\left\{k \in Y: p H \leq s_{k}<\left(p+\frac{1}{3}\right) H\right\} \tag{2.23}
\end{equation*}
$$

(It is eq. (2.22) that depends crucially on the support properties.) Let

$$
\begin{equation*}
\Lambda_{j}(p)=\sum_{k \in Y_{p}} \lambda_{j}(k) \tag{2.24}
\end{equation*}
$$

Since $P_{*}$ is norm decreasing, eqs. (2.12), (2.20) and (2.21)-(2.24) give

$$
\begin{equation*}
\left\|\mu_{j}-\sum_{p \in \mathbb{Z}} \Lambda_{j}(p) T^{p} \sigma\right\|<8 H^{-1} \varepsilon \tag{2.25}
\end{equation*}
$$

where $\sigma=6 H^{-1} P_{*} \mu^{\prime}$.
The tower construction of a single transformation (see for example [9]) is the analogue of the base and ceiling function construction of a flow.
Corollary 2.6. Let the transformation ( $X, \nu, S$ ) be constructed as a tower over ( $B, \nu_{B}, T$ ) with constant height $H$. Then $S$ is AT if and only if $T$ is AT.
Proof. The flow $F_{s}$ built over ( $X, \nu, S$ ) with constant height one is obviously a flow built over ( $B, \nu_{B}, T$ ) with constant height $H$. The result now follows from lemma 2.5.
Recall that a finite measure preserving transformation is said to have rank one if there is a sequence of Rohlin towers which approximate the measure algebra. More precisely one asks that given $f_{1}, \ldots, f_{n} \in L^{\prime}(x, \nu)$ and $\varepsilon>0$, there exist $B \subset X, m<\infty$, and $\lambda_{j k} \in \mathbb{R}$ such that $B, T B, \ldots, T^{m} B$ are disjoint and

$$
\begin{equation*}
\left\|f_{j}-\sum_{k=0}^{m} \lambda_{j k} X_{T^{k} B}\right\|<\varepsilon, \quad j=1, \ldots, n . \tag{2.26}
\end{equation*}
$$

(Note that if $f_{j} \geq 0$ we can require that $\lambda_{j k} \geq 0$.) A slightly weaker condition, called funny rank one, is obtained by replacing the sequence $B, T B, \ldots, T^{m} B$ by the sets $T^{n_{0}} B, T^{n_{1}} B, \ldots, T^{n_{m}} B$ where $\left\{n_{j}\right\}$ is an arbitrary sequence (depending on $f_{1}, \ldots, f_{n}$ and $\varepsilon\}$.

Lemma 2.7. Given $(X, T, \nu)$ with $T \nu=\nu, \nu(X)=1$. Then funny rank one implies AT. Proof. It is convenient here to use the AT condition on functions in $L_{+}^{1}(X, \nu)$ (see remark 2.2 (iv)). Let $f_{1}, \ldots, f_{n} \in L_{+}^{1}(X, \nu), \varepsilon>0$. Then there exist $B \subset X$, a sequence $\left\{n_{k}\right\}_{k=0, \ldots, m}$ and $\lambda_{j k}$ such that

$$
\begin{equation*}
\left\|f_{j}-\sum_{k=1}^{m} \lambda_{j k} X T_{T^{n} \times B}\right\|>\varepsilon . \tag{2.27}
\end{equation*}
$$

Let

$$
\lambda_{j k}^{\prime}= \begin{cases}\lambda_{j k} & \text { if } \lambda_{j k} \geq 0  \tag{2.28}\\ 0 & \text { if } \lambda_{j k}<0\end{cases}
$$

Since the $T^{n_{k}} B$ are disjoint, we have

$$
\begin{equation*}
\left\|f_{j}-\sum \lambda_{j k}^{\prime} T^{-n_{k}} f\right\| \leq\left\|f_{j}-\sum \lambda_{j k} T^{-n_{k}} f\right\|, \tag{2.29}
\end{equation*}
$$

where $f=\chi_{B}$.
Corollary 2.8. If $(X, T, \nu)$ has pure point spectrum then $T$ is AT.
Proof. Pure point spectrum implies rank one [11].
Corollary 2.9. A pure point spectrum flow is AT.
Proof. Such a flow can be built over a pure point base transformation with a constant ceiling function. The result now folllows from corollary 2.8 and lemma 2.5.
It is also known that certain diffeomorphisms of the circle are AT [10].

## 3. AT transformations and entropy

In this section we prove that if $(X, \mu, T)$ is a finite measure preserving transformation, then AT implies that $T$ has zero entropy (theorem 3.5).

In analyzing the implications of the AT condition one immediately observes that an expression of the form

$$
\begin{equation*}
T(\lambda) f=\sum \lambda_{j} T^{j} f, \tag{3.1}
\end{equation*}
$$

where $\lambda \in \ell_{+}^{\prime}(\mathbb{Z})$ and $f \in L_{+}^{1}(X, \mu)$, is in effect a convolution, and 'convolutions spread functions out'. This spreading can present some difficulties when one tries to satisfy the AT condition for functions $f_{1}=\chi_{A_{1}}, f_{2}=\chi_{A_{2}}$ where $A_{1}, A_{2} \subset X$. In order to give a precise meaning to the idea that $T(\lambda) f$ is 'less concentrated' than $f$, we define upper and lower truncations of $L^{1}$ functions (definition 3.1). The upper truncation is used to measure the concentration of $f$ (see eq. (3.5)), and lemma 3.2 then gives a precise meaning to the statement that convolutions spread. We give a 'spectral analysis' of $L^{1}$ functions, which allows one to handle the difficulties that arise in an argument when $L^{1}$ functions take on either very large or very small values.

Lemmas 3.3 and 3.4 are technical lemmas required for the proof of theorem 3.5. (Hint: read the proof of theorem 3.5 before reading lemmas 3.3 and 3.4.) The basic idea of the argument is as follows. One applies the AT condition to $f_{1}=\chi_{\mathrm{A}}$ for some $A \subset X$, a second function $f_{2}$, and some $\varepsilon>0$. This gives functions $\lambda_{1}, \lambda_{2}$ and $f$ satisfying $\left\|T\left(\lambda_{i}\right) f-f_{i}\right\| \leq \varepsilon, i=1,2$. Since convolutions spread, choosing $f_{2}$ very concentrated relative to the set $A$ forces $f$ to be very concentrated relative to $f_{1}$.

This in turn forces $\lambda_{1}$ to be very 'spread out' (lemma 3.3). However $\lambda_{1}$ being very spread out makes it difficult to keep the support of $T\left(\lambda_{1}\right) f$ close to $A$, which it must approximate, and simultaneously to keep $T\left(\lambda_{1}\right) f$ small on $A^{c}$. In particular, when the partition $(A, X \backslash A$ ) moves independently under $T$, this becomes impossible (theorem 3.5). In order to make this argument precise, one must replace $\lambda_{1}$ and $f$ by functions $\lambda \cong \lambda_{1}$ and $g \cong f$ with better support properties. This is done in lemma 3.4 by a technical application of the spectral analysis for $L^{1}$ functions. The condition that $T(\lambda) g$ be small on $A^{c}$ then forces the support of $T(\lambda) g$ to be too small to approximate $A$. In particular, the proof seems somewhat stronger than the statement of the theorem. It suffices that $k \mu(B) \rightarrow 0$ (see eqs. (3.30), (3.31), (3.41), (3.42), (3.52), (3.53)).

Let $(X, \mu)$ be a Lebesgue measure space. Then $(X, \mu)$ is isomorphic to Lebesgue measure on $[0,1]$. Let $f \in L_{+}^{1}(X, \mu)$. Then one can choose the isomorphism so that the (transformed) function $f$ is monotone decreasing. The upper (resp. lower) truncation of $f$ is a function whose graph is identical to the graph of $f$, except that the upper left hand corner (resp. lower right hand corner) has been 'chopped off'. More precisely we have:
Definition 3.1. Let $f \in L_{+}^{1}(X, \mu), a>0$. We define the upper truncation of $f$ at $a$ by

$$
f^{[a]}(x)= \begin{cases}f(x) & \text { if } f(x) \leq a  \tag{3.2}\\ a & \text { if } f(x)>a\end{cases}
$$

and the lower truncation of $f$ at $a$ by

$$
f_{[a]}(x)= \begin{cases}f(x) & \text { if } f(x) \geq a  \tag{3.3}\\ 0 & \text { if } f(x)<a\end{cases}
$$

The continuity of upper truncations is expressed by the condition

$$
\begin{equation*}
\|f-g\| \leq \varepsilon \Rightarrow\left\|f^{[a]}-g^{[a]}\right\| \leq \varepsilon \tag{3.4}
\end{equation*}
$$

which follows immediately from eq. (3.2). It should be noted that eq. (3.4) does not hold for lower truncations.

Consider the inequality

$$
\begin{equation*}
\left\|f-f^{[a]}\right\| \geq\|f\|-\eta \tag{3.5}
\end{equation*}
$$

where $a, \eta>0$. If $\eta$ is small compared to $\|f\|$, this inequality forces most of the contribution to the $L^{1}$ norm of $f$ to come from that part of the graph of $f$ lying above $a$. If in addition $a$ is large compared to $\|f\|$, it then forces most of $f$ (in the sense of $L^{1}$ norm) to be supported on a set of small measure. It can therefore be used as a measure of the 'concentration' of a function. The following lemma now gives a precise meaning to the statement that 'convolutions spread'.
Lemma 3.2. Given $(X, \mu, T), \lambda \in \ell_{+}^{1}(\mathbb{Z}),\|\lambda\|=1$, and $f \in L_{+}^{1}(X, \mu)$, then

$$
\begin{equation*}
\| T(\lambda) f-(T \lambda) f)^{[a]}\|\leq\| f-f^{[a]} \| \tag{3.6}
\end{equation*}
$$

Proof. If $f_{j}, \sum \lambda_{j} f_{j} \in L_{+}^{1}(X, \mu)$ then it follows directly from eq. (3.2) that

$$
\begin{equation*}
\left(\sum \lambda_{j} f_{j}\right)^{[a]}(x) \geq \sum \lambda_{j}\left(f_{j}^{[a]}\right)(x) . \tag{3.7}
\end{equation*}
$$

Since $\left\|f-f^{[a]}\right\|=\|f\|-\left\|f^{[a]}\right\|$, and $\|T(\lambda) f\|=\|f\|$, eq. (3.6) follows immediately.

We now proceed to the spectral analysis of $L^{1}$ functions. This analysis is, in effect, a consequence of considering $\int f d \mu$ in terms of horizontal slices of the graph of $f$. Let $f \in L_{+}^{1}(X, \mu)$. For each $a>0$ we define a function $E_{a}(f)$ on $X$ by

$$
\left(E_{a}(f)\right)(x)= \begin{cases}1 & \text { if } f(x) \geq a .  \tag{3.8}\\ 0 & \text { if } f(x)<a .\end{cases}
$$

We define a measure $\nu_{f}$ on $\mathbb{R}$, absolutely continuous with respect to Lebsgue measure $d a$, by

$$
\begin{equation*}
\left(d \nu_{f} / d a\right)(a)=\int E_{a}(f) d \mu \tag{3.9}
\end{equation*}
$$

It follows that a.e. we have

$$
\begin{equation*}
f(x)=\int_{0}^{\infty}\left(E_{a}(f)\right)(x) d a=\int_{0}^{\infty} E_{a}^{\prime}(f) d \nu_{f}(a) \tag{3.10}
\end{equation*}
$$

where

$$
E_{a}^{\prime}(f)= \begin{cases}0 & \text { if } \int E_{a}(f) d \mu=0  \tag{3.11}\\ E_{a}(f)\left[\int E_{a}(f) d \mu\right]^{-1} & \text { otherwise }\end{cases}
$$

We then have

$$
\begin{gather*}
\|f\|=\int f d \mu=\int d \nu_{f} \int d \mu E_{a}^{\prime}(f)=\int d \nu_{f}  \tag{3.12}\\
\left\|f-f^{[a]}\right\|=\int_{a}^{\infty} d \nu_{f}\left(a^{\prime}\right) \tag{3.13}
\end{gather*}
$$

and

$$
\begin{equation*}
\left\|f-f_{[a]}\right\| \leq \int_{0}^{a} d \nu_{f}\left(a^{\prime}\right) \tag{3.14}
\end{equation*}
$$

(Equality in eq. (3.14) holds only when $f=0$.) These equations indicate the significance of the measure $\nu_{f}$.
Lemma 3.3. Given $(X, \mu, T), \mu \circ T=\mu, f_{1} \in L_{+}^{1}(X, \mu),\left\|f_{1}\right\| \leq 1$ and $\varepsilon>0$, then there exists $f_{2} \in L_{+}^{1}(X, \mu),\left\|f_{2}\right\|=1$, such that for any $f \in L_{+}^{1}(X, \mu),\|f\|=1$, and $\lambda_{1}, \lambda_{2} \in$ $\ell_{+}^{1}(\mathbb{Z}),\left\|\lambda_{1}\right\| \leq 1,\left\|\lambda_{2}\right\|=1$ satisfying

$$
\begin{equation*}
\left\|T\left(\lambda_{i}\right) f-f_{i}\right\| \leq \varepsilon, \quad i=1,2 \tag{3.15}
\end{equation*}
$$

we have

$$
\begin{equation*}
\sup _{j \in Z} \lambda_{1 j} \leq 6 \varepsilon . \tag{3.16}
\end{equation*}
$$

Proof. The proof of this lemma consists of a long sequence of inequalities. Nevertheless it is a completely straightforward and obvious application of the above ideas. Choose $a<\infty$ such that

$$
\begin{equation*}
\left\|f_{1}-f_{1}^{[a]}\right\| \leq \varepsilon \tag{3.17}
\end{equation*}
$$

Choose $f_{2} \in L_{+}^{1}(X, \mu),\left\|f_{2}\right\|=1$, such that

$$
\begin{equation*}
\left\|f_{2}-f_{2}^{\left[a \varepsilon^{-1}\right]}\right\| \geq 1-\varepsilon \tag{3.18}
\end{equation*}
$$

It follows from eqs. (3.4) and (3.15) that

$$
\begin{equation*}
\left\|f_{2}^{\left[a \varepsilon^{-1}\right]}-\left(T\left(\lambda_{2}\right)\right) f^{\left[a \varepsilon^{-1}\right]}\right\| \leq \varepsilon \tag{3.19}
\end{equation*}
$$

Eqs. (3.15), (3.18) and (3.19) give

$$
\begin{equation*}
\left\|T\left(\lambda_{2}\right) f-\left(T\left(\lambda_{2}\right) f\right)^{\left[a \varepsilon^{-1}\right]}\right\| \geq 1-3 \varepsilon \tag{3.20}
\end{equation*}
$$

Eq. (3.20) and lemma 3.2 give

$$
\begin{equation*}
\left\|f-f^{\left[a \varepsilon^{-1}\right]}\right\| \geq 1-3 \varepsilon \tag{3.21}
\end{equation*}
$$

Let

$$
\begin{equation*}
A=\left\{x: f(x) \geq a \varepsilon^{-1}\right\} . \tag{3.22}
\end{equation*}
$$

Since $\|f\| \leq 1$ we have

$$
\begin{equation*}
\mu(A) \leq a^{-1} \varepsilon \tag{3.23}
\end{equation*}
$$

Eq. (3.15) gives for each $j \in \mathbb{Z}$,

$$
\begin{align*}
\varepsilon \geq\left\|f_{1}-\sum \lambda_{1 j} T^{j} f\right\| & \geq \int_{T^{-j} A}\left(\sum \lambda_{1 j} T^{j} f-f_{1}\right) d \mu \\
& \geq \int_{T^{-j} A}\left(\lambda_{1 j} T^{j} f-f_{1}\right) d \mu \tag{3.24}
\end{align*}
$$

Eqs. (3.21) and (3.22) give

$$
\begin{equation*}
\int_{T^{-j} A} \lambda_{1 j} T^{j} f d \mu \geq \lambda_{1 j}(1-3 \varepsilon) \tag{3.25}
\end{equation*}
$$

Eqs. (3.17) and (3.23) give

$$
\begin{equation*}
\int_{T^{-j_{A}}} f_{1} d \mu \leq 2 \varepsilon \tag{3.26}
\end{equation*}
$$

Eqs. (3.24)-(3.26) give $(1-3 \varepsilon) \lambda_{1 j} \leq 3 \varepsilon$ and hence

$$
\lambda_{1 j} \leq 3 \varepsilon\left(1+\lambda_{1 j}\right) \leq 6 \varepsilon .
$$

Since convolutions preserve $L^{1}$ norms and $\|f\|=1$, it follows from eq. (3.15) that $\left\|\lambda_{1}\right\| \geq\left\|f_{1}\right\|-\varepsilon$. Thus if $\varepsilon$ is very small, eq. (3.16) forces the support of $\lambda_{1}$ to be very large.
Lemma 3.4. Given $(X, \mu, T), A \subset X, f \in L_{+}^{1}(X, \mu), \lambda_{1} \in \ell_{+}^{1}(\mathbb{Z})$ and $\varepsilon>0$ such that

$$
\begin{equation*}
\left\|T\left(\lambda_{1}\right) f-\chi_{A}\right\| \leq \varepsilon \tag{3.27}
\end{equation*}
$$

then there exist $g \in L_{+}^{1}(X, \mu), \lambda \in \ell_{+}^{1}(\mathbb{Z})$ such that

$$
\begin{align*}
\|f-g\| & \leq 4 \varepsilon^{\frac{1}{2}},  \tag{3.28}\\
\left\|\lambda_{1}-\lambda\right\| & \leq \varepsilon^{\frac{1}{2}}, \tag{3.29}
\end{align*}
$$

and

$$
\begin{equation*}
\mu(S) \leq k \mu(B) \tag{3.30}
\end{equation*}
$$

where $S$ is the support of $T(\lambda) g, k$ is the number of elements in the support $K$ of $\lambda$, and

$$
\begin{equation*}
B=\left\{x \in X: k^{-1} \sum_{j \in K} \chi_{A^{c}}\left(T^{j} x\right)<\frac{1}{4}\right\} . \tag{3.31}
\end{equation*}
$$

Proof. Eq. (3.27) implies that

$$
\begin{equation*}
\int_{A^{c}} T\left(\lambda_{\mathrm{I}}\right) f d \mu \leq \varepsilon \tag{3.32}
\end{equation*}
$$

It follows from eq. (3.10) that

$$
\begin{equation*}
T\left(\lambda_{1}\right)=\int d \nu_{\lambda_{1}}(a) T\left(E_{a}^{\prime}\left(\lambda_{1}\right)\right) \tag{3.33}
\end{equation*}
$$

Let

$$
\begin{equation*}
\rho(a)=\int_{A^{c}} T\left(E_{a}^{\prime}(\lambda)\right) f d \mu \tag{3.34}
\end{equation*}
$$

Eqs. (3.32)-(3.34) give

$$
\begin{equation*}
\int d \nu_{\lambda_{1}}(a) \rho(a) \leq \varepsilon \tag{3.35}
\end{equation*}
$$

and hence

$$
\begin{equation*}
\nu_{\lambda_{1}}\left\{a: \rho(a)>\varepsilon^{\frac{1}{2}}\right\}<\varepsilon^{\frac{1}{2}} . \tag{3.36}
\end{equation*}
$$

It follows from eqs. (3.34)-(3.36) that there exists $b>0$ such that

$$
\begin{equation*}
\rho(b) \leq \varepsilon^{\frac{1}{2}} \tag{3.37}
\end{equation*}
$$

and

$$
\begin{equation*}
\int_{0}^{b} d \nu_{\lambda_{1}}(a) \leq \varepsilon^{\frac{1}{2}} . \tag{3.38}
\end{equation*}
$$

It follows that $\lambda=\left(\lambda_{1}\right)_{[b]}$ satisfies eq. (3.29). Since $K$ is also the support of $E_{b}^{\prime}(\lambda)$, eq. (3.37) becomes

$$
\begin{align*}
\varepsilon^{\frac{1}{2}} & \geq \int_{A^{c}} d \mu k^{-1} \sum_{j \in K} T^{j} f=k^{-1} \sum_{j \in K} \int_{T^{-j} A^{c}} d \mu f \\
& =\int_{X} d \mu\left[k^{-1} \sum_{j \in K} \chi_{A^{c}}\left(T^{j} x\right)\right] f(x) . \tag{3.39}
\end{align*}
$$

Eqs. (3.31) and (3.39) give

$$
\begin{equation*}
\int_{B^{c}} d \mu f \leq 4 \varepsilon^{\frac{1}{2}} . \tag{3.40}
\end{equation*}
$$

It follows that $g(x)=\chi_{B}(x) f(x)$ satisfies eq. (3.28). Eq. (3.30) is satisfied by construction.

Theorem 3.5. Let $(X, \mu, T)$ be an AT transformation, $\mu(X)=1$, and $\mu \circ T=\mu$. Then the entropy $h(T)=0$.
Proof. Assume that $h(T)>0$. By a well-known result of Sinai there is a partition $\left\{A, A^{c}\right\}$ of $X$ that moves independently under $T$ (see for example [18, p. 43]). We can assume that $0<\mu(A) \leq \frac{1}{2}$. Consider the set $B$ given by lemma 3.4 (where $\varepsilon, f$,
and $\lambda_{1}$ will be chosen below). Then $\mu(B)$ depends only on the number $k$ of elements in $K$, and can be calculated directly from the binomial coefficients. Since the special case $\mu(A)=\frac{1}{2}$ dominates, one easily obtains the inequality

$$
\begin{equation*}
\mu(B)<2^{-k} . \tag{3.41}
\end{equation*}
$$

Choose $k_{0}$ sufficiently large that

$$
\begin{equation*}
k_{0} 2^{-k_{0}} \leq \frac{1}{2} \mu(A) \tag{3.42}
\end{equation*}
$$

and choose $\varepsilon>0$ sufficiently small that

$$
\begin{equation*}
\varepsilon+5 \varepsilon^{\frac{1}{2}} \leq \frac{1}{2} \mu(A) \tag{3.43}
\end{equation*}
$$

and

$$
\begin{equation*}
12 \varepsilon k_{0} \leq \mu(A) \tag{3.44}
\end{equation*}
$$

Now let $f_{1}(x)=\chi_{A}(x)$. By lemma 3.3 there exists $f_{2} \in L_{+}^{1}(X, \mu),\left\|f_{2}\right\|=1$, such that eq. (3.16) is satisfied. Since $T$ is AT there exists $f \in L_{+}^{1}(X, \mu),\|f\|=1$ and $\lambda_{1}, \lambda_{2} \in$ $\ell_{+}^{1}(\mathbb{Z}),\left\|\lambda_{2}\right\|=1$, such that

$$
\begin{equation*}
\left\|T\left(\lambda_{i}\right) f-f_{i}\right\| \leq \varepsilon, \quad i=1,2 . \tag{3.45}
\end{equation*}
$$

By lemma 3.4 there exist $\lambda$ and $g=\chi_{B} f$ satisfying eqs. (3.28)-(3.30). Eqs. (3.28)(3.29) and (3.43) give

$$
\begin{equation*}
\left\|T(\lambda) g-\chi_{A}\right\| \leq \frac{1}{2} \mu(A) \tag{3.46}
\end{equation*}
$$

and hence

$$
\begin{equation*}
\mu(S) \geq \frac{1}{2} \mu(A) \tag{3.47}
\end{equation*}
$$

where $S$ is the sup;ort of $T(\lambda) g$. Since $\|g\| \leq\|f\|=1$ and convolutions preserve $L^{1}$ norms, it also follows from eq. (3.46) that

$$
\begin{equation*}
\|\lambda\| \geq \frac{1}{2} \mu(A) \tag{3.48}
\end{equation*}
$$

Since $\lambda=\left(\lambda_{1}\right)_{[b]}$, eq. (3.16) implies that we also have

$$
\begin{equation*}
\sup _{j \in \mathbb{Z}} \lambda(j) \leq 6 \varepsilon . \tag{3.49}
\end{equation*}
$$

Eqs. (3.48)-(3.49) give

$$
\begin{equation*}
6 \varepsilon k \geq \frac{1}{2} \mu(A) \tag{3.50}
\end{equation*}
$$

where $k$ is the number of elements in the support $K$ of $\lambda$. Eqs. (3.44) and (3.50) give

$$
\begin{equation*}
k \geq k_{0} . \tag{3.51}
\end{equation*}
$$

Eqs. (3.51), (3.41) and (3.42) give

$$
\begin{equation*}
k \mu(B)<\frac{1}{2} \mu(A) \tag{3.52}
\end{equation*}
$$

Eq. (3.30), which was obtained from the requirement that $T\left(\lambda_{1}\right) f$ is small on $A^{c}$, now gives

$$
\begin{equation*}
\mu(S)<\frac{1}{2} \mu(A) \tag{3.53}
\end{equation*}
$$

which contradicts eq. (3.47).

## 4. Constructions of the flow of weights

In order to make our exposition reasonably self-contained, we give here three different constructions of the flow of weights, each of which will be used at some
point in our argument. We also construct some measures on the resulting spaces. For the proofs of all statements in this section, see [5].
Discrete construction. Let $M$ be a type $\mathrm{III}_{0}$ factor acting on the Hilbert space $H$, with $T \in T(M), T>0$. (If $T(M)=\{0\}$ one is then forced to use one of the following two constructions.) Let $\phi$ be a faithful state on $M, \sigma_{T}^{\phi}=1$. Let $\left(e_{n}\right)_{n \in \mathbb{Z}}$ be the canonical orthonormal basis for $\ell^{2}(\mathbb{Z})$, and define

$$
\begin{align*}
S e_{n} & =e_{n+1}  \tag{4.1}\\
\rho_{\lambda} e_{n} & =\lambda^{n} e_{n} \tag{4.2}
\end{align*}
$$

where $n \in \mathbb{Z}$ and $\lambda=\exp (-2 \pi / T)$. The equation $\omega_{\lambda}(x)=$ Trace $\rho_{\lambda} x$ defines a faithful semifinite normal weight $\omega_{\lambda}$ on $\mathscr{L}\left(\ell^{2}(\mathbb{Z})\right)$. We have

$$
\begin{equation*}
\sigma_{\lambda}^{\omega_{\lambda}}(S)=\rho_{\lambda}^{i t} S \rho_{\lambda}^{-i t}=\lambda^{i t} S, \tag{4.3}
\end{equation*}
$$

and

$$
\begin{equation*}
\omega_{\lambda}\left(S A S^{*}\right)=\lambda \omega_{\lambda}(A) . \tag{4.4}
\end{equation*}
$$

Let

$$
\begin{align*}
\tilde{H} & =\ell^{2}(\mathbb{Z}) \otimes H,  \tag{4.5}\\
\tilde{M} & =\mathscr{L}\left(\ell^{2}(\mathbb{Z})\right) \otimes M,  \tag{4.6}\\
\tilde{\phi} & =\omega_{\lambda} \otimes \phi,  \tag{4.7}\\
\tilde{S} & =S \otimes 1 . \tag{4.8}
\end{align*}
$$

Then $\sigma_{t}^{\dot{\phi}}=\sigma_{1}^{\omega_{\lambda}} \otimes \sigma_{t}^{\phi}$ and it follows from eq. (4.3) that the automorphism $\theta=\operatorname{Ad} \tilde{S}$ leaves the centralizer $N=\tilde{M}_{\tilde{\Phi}}$ invariant. Eqs. (4.4), (4.7) and (4.8) give

$$
\begin{equation*}
\tilde{\phi} \circ \theta=\lambda \tilde{\phi} . \tag{4.9}
\end{equation*}
$$

The centre $C$ of $N$ is isomorphic to $L^{\infty}\left(B, \nu_{B}\right)$ where $\left(B, \nu_{B}\right)$ is a Lebesgue measure space. $\theta$ then defines an automorphism, which we shall also denote by $\theta$, of $\left(B, \nu_{B}\right)$. The flow of weights for $M$ is the flow ( $X, \nu, F_{t}$ ) built over the base transformation ( $B, \nu_{B}, \theta$ ) with a ceiling function of constant height $2 \pi / T$. Furthermore

$$
\begin{equation*}
N=\int_{B}^{\oplus} N(b) d \nu_{B}(b), \tag{4.10}
\end{equation*}
$$

where the $N(b)$ are type $\mathrm{II}_{\infty}$ factors. If $M$ is injective then the $N(b)$ are all (isomorphic to) the unique injective $\mathrm{II}_{\infty}$ factor $R_{0,1}$ (see [4]). One can then write

$$
\begin{equation*}
N=R_{0,1} \otimes L^{\infty}\left(B, \nu_{B}\right)=\int_{B}^{\oplus} R_{0,1} d \nu_{B}(b), \tag{4.11}
\end{equation*}
$$

and

$$
\begin{equation*}
\tilde{\phi}=\tau \otimes \nu_{B}=\int^{\oplus} \tau_{b} d \nu_{B}(b), \tag{4.12}
\end{equation*}
$$

where $\tau$ is a trace on $R_{0,1}$ and $\tau_{b}$ is a trace on $N(b)$.
We now assign to certain positive operators in $N$, measures $\mu$ on $B, \mu<\nu_{B}$. Let

$$
\begin{equation*}
A=\int_{B}^{\oplus} A(b) d \nu_{B}(b) \in N_{+} \tag{4.13}
\end{equation*}
$$

be such that

$$
\begin{equation*}
\tilde{\phi}(A)=\int_{B} \tau_{b}(A(b)) d \nu_{B}(b)<\infty \tag{4.14}
\end{equation*}
$$

The equation

$$
\begin{equation*}
\mu_{A}(c)=\tilde{\phi}(A c), \quad c \in C \tag{4.15}
\end{equation*}
$$

defines a finite measure $\mu_{A}<\nu_{B}$. It follows from eqs. (4.14) and (4.15) that

$$
\begin{equation*}
\left(d \mu_{A} / d \nu_{B}\right)(b)=\tau_{b}(A(b)) \tag{4.16}
\end{equation*}
$$

It follows from eq. (4.16) that if, $e, f$ are $\tilde{\phi}$-finite projections in $N$, then $\mu_{e}=\mu_{f}$ if and only if $e \sim f$ in $N$. Since $N(b) \sim N(b) \otimes R_{0,1}$ one can construct families of projections $e_{b}(\alpha) \in N(b), \alpha>0$ such that

$$
\begin{equation*}
\tau_{b}\left(e_{b}(\alpha)\right)=\alpha, \tag{4.17}
\end{equation*}
$$

and $e_{b}(f(b))$ is a $\nu_{B}$-measurable family of projections for any $f \in L_{+}^{1}\left(B, \nu_{B}\right)$. It follows that every finite measure $\mu<\nu_{B}$ occurs as $\mu_{e}$ for some $e$, namely

$$
\begin{equation*}
e=\int_{B}^{\oplus} e_{b}\left(d \mu / d \nu_{B}\right) d \nu_{B} . \tag{4.18}
\end{equation*}
$$

Continuous construction. Let $M$ be a type $\mathrm{III}_{0}$ factor acting on the Hilbert space $H$, $\phi$ a faithful state on $M$. On $L^{2}(\mathbb{R})$ we define

$$
\begin{equation*}
\left(V_{s} f\right)(t)=f(t-s), \tag{4.19}
\end{equation*}
$$

and

$$
\begin{equation*}
(\rho f)(t)=e^{t} f(t) \tag{4.20}
\end{equation*}
$$

The equation $\omega(x)=$ Trace $\rho x$ defines a faithful semifinite normal weight $\omega$ on $\mathscr{L}\left(L^{2}(\mathbb{R})\right)$. We have

$$
\begin{equation*}
\sigma_{t}^{\omega}\left(V_{s}\right)=e^{-i t s} V_{s}, \tag{4.21}
\end{equation*}
$$

and

$$
\begin{equation*}
\omega\left(V_{s} A V_{s}^{*}\right)=e^{t} \omega(A) \tag{4.22}
\end{equation*}
$$

Let

$$
\begin{align*}
\tilde{H} & =L^{2}(\mathbb{R}) \otimes H  \tag{4.23}\\
\tilde{M} & =\mathscr{L}\left(L^{2}(\mathbb{R})\right) \otimes M,  \tag{4.24}\\
\tilde{\omega} & =\omega \otimes \phi  \tag{4.25}\\
W_{s} & =V_{s} \otimes 1 \tag{4.26}
\end{align*}
$$

Then $\theta_{s} N \theta_{s}^{-1}=N$ where $N=\tilde{M}_{\tilde{\omega}}$ and $\theta_{s}=\operatorname{Ad} W_{s}$. We have

$$
\begin{equation*}
N=\int_{X}^{\oplus} N(x) d \nu(x), \tag{4.27}
\end{equation*}
$$

where the $N(x)$ are type $\mathrm{I}_{\infty}$ factors, and the centre $C$ of $N$ is isomorphic to $L^{\infty}(X, \nu)$ where $(X, \nu)$ is a Lebesgue measure space. The automorphisms $\theta_{s}$ define a flow $F_{s}$ on $(X, \nu)$ which is the flow of weights for $M$. If $M$ is injective then
$N(x) \sim R_{0,1}[4]$, and one can write

$$
\begin{equation*}
N=R_{0,1} \otimes L^{\infty}(X, \nu)=\int_{X}^{\oplus} R_{0,1} d \nu(x) \tag{4.28}
\end{equation*}
$$

and

$$
\begin{equation*}
\tilde{\omega}=\tau \otimes \nu=\int^{\oplus} \tau d \nu \tag{4.29}
\end{equation*}
$$

where $\tau$ is a trace on $R_{0,1}$.
The construction of measures in the continuous construction is quite analogous to the discrete case, except that they occur on the flow space $X$ rather than the base space $B$. Let

$$
\begin{equation*}
A=\int_{X}^{\oplus} A(x) d \nu(x) \in N_{+} \tag{4.30}
\end{equation*}
$$

be such that

$$
\begin{equation*}
\tilde{\omega}(A)=\int_{X} \tau_{x}(A(x)) d \nu(x)<\infty \tag{4.31}
\end{equation*}
$$

Then the equation

$$
\begin{equation*}
\mu_{A}(c)=\tilde{\omega}(A c), \quad c \in C \tag{4.32}
\end{equation*}
$$

defines a finite measure $\mu_{A}<\nu$ such that

$$
\begin{equation*}
\left(d \mu_{A} / d \nu\right)(x)=\tau_{b}(A(x)) \tag{4.33}
\end{equation*}
$$

If $e, f$ are $\tilde{\omega}$-finite projections in $N$, then $\mu_{e}=\mu_{f}$ if and only if $e \sim f$ in $N$. As in the discrete case, every finite measure $\mu<\nu$ occurs as $\mu_{e}$ for some projection $e$. We shall need the following lemma.

Lemma 4.1. Let $(Y, \sigma)$ be a measure space, and let $e_{y}$ be a $\sigma$-measurable family of $\tilde{\omega}$-finite positive operators in $N$ such that

$$
\begin{equation*}
e=\int_{Y} e_{y} d \sigma(y) \tag{4.34}
\end{equation*}
$$

is $\tilde{\omega}$-finite. Then

$$
\begin{equation*}
\mu_{e}=\int_{Y} \mu_{e_{y}} d \sigma(y) \tag{4.35}
\end{equation*}
$$

Proof. We have

$$
\begin{equation*}
e=\int_{X}^{\oplus} e(x) d \nu(x) \tag{4.36}
\end{equation*}
$$

and

$$
\begin{equation*}
e_{y}=\int_{X}^{\oplus} e_{y} d \nu(x) \tag{4.37}
\end{equation*}
$$

Eqs. (4.34), (4.36) and (4.37) imply that

$$
\begin{equation*}
e(x)=\int_{Y} e_{y}(x) d \sigma(y) \quad(\text { a.e. } \nu) \tag{4.38}
\end{equation*}
$$

Eq. (4.38) implies eq. (4.35).

Lacunary construction. Let $M$ be a factor of type $\mathrm{III}_{0}, \phi$ a faithful lacunary weight of infinite multiplicity on $M$ (lacunary means that 1 is an isolated point in $\operatorname{Sp} \Delta_{\phi}$ ). Then

$$
\begin{equation*}
M_{\phi}=\int_{B}^{\oplus} M(b) d \nu_{B}(b), \tag{4.39}
\end{equation*}
$$

where the $M(b)$ are type $\mathrm{I}_{\infty}$ factors, and the centre $C_{\phi}$ of the centralizer $M_{\phi}$ is isomorphic to $L^{\infty}\left(B, \nu_{B}\right)$ where $\left(B, \nu_{B}\right)$ is a Lebesgue measure space. If $M$ is injective then $M(b) \sim R_{0,1}$. There exists $\rho \in C_{\phi}, 0<\rho \leq \lambda_{0}$ for some $\lambda_{0}<1$, and a unitary $U \in M$ such that

$$
\begin{align*}
\phi\left(U x U^{*}\right) & =\phi(\rho x), \quad x \in M,  \tag{4.40}\\
M & =\left\{M_{\phi}, U\right\}^{\prime \prime}, \tag{4.41}
\end{align*}
$$

and

$$
\begin{equation*}
U M_{\phi} U^{*}=M_{\phi} . \tag{4.42}
\end{equation*}
$$

Then $\theta=\operatorname{Ad} U$ defines an automorphism, which we also denote by $\theta$, of $\left(B, \nu_{B}\right)$. The flow of weights for $M$ is the flow ( $X, \nu, \mathscr{F}_{s}, s \in \mathbb{R}_{+}^{*}$ ) built over $\left(B, \nu_{B}, \theta\right.$ ) with the ceiling function $\rho$, where

$$
\begin{equation*}
X=\{(b, t): b \in B, 1 \geq t>\rho(b)\} \tag{4.43}
\end{equation*}
$$

and $\mathscr{F}_{s}(b, t)=\left(b, e^{-s} t\right)$ if $1 \geq e^{-s} t>\rho(b)$ with the obvious extension to other values of $s$.

We again construct certain measures on the flow space $X$. Let $\psi \in \boldsymbol{M}_{*}^{+}$. Then there exists $h \in M_{\phi}^{+}$such that $s(h) \rho \leq h \leq 1,1-h$ is non-singular, and there exists a unitary $u \in M$ such that

$$
\begin{equation*}
\psi(x)=\phi\left(h E\left(u x u^{*}\right)\right) \tag{4.44}
\end{equation*}
$$

where $E$ is the conditional expectation from $M$ onto $M_{\phi}$. We have

$$
\begin{equation*}
h=\int_{B}^{\oplus} h_{b} d \nu_{B}(b) \tag{4.45}
\end{equation*}
$$

where $s\left(h_{b}\right) \rho(b) \leq h_{b} \leq 1$. Let $f \in L^{\infty}(X, \nu)$. Then the operator

$$
\begin{equation*}
h f(h)=\int_{B}^{\oplus} h_{b} f_{b}\left(h_{b}\right) d \nu_{B}(b), \tag{4.46}
\end{equation*}
$$

where $f_{b}(t)=f(b, t)$, is well-defined. The equation

$$
\begin{equation*}
\mu_{\psi}(f)=\phi(h f(h)) \tag{4.47}
\end{equation*}
$$

defines a measure $\mu_{\psi}$ on the flow space $X$ (which we will sometimes write as $\mu_{h}$ ). (Eqs. (4.47) and (4.32) are related as follows. Let $\psi=\tilde{\omega}_{e}, e \in \tilde{M}_{\tilde{\omega}}$. Then $\mu_{\psi}=\mu_{e}$ ) Some terminology is helpful at this point.
Definition 4.2. Let $\left(X, \nu, F_{t}\right)$ be a flow. We call the measure $\mu$ on $X$ smooth if $\mu<\nu$, and smoothable if

$$
\begin{equation*}
f * \mu=\int_{-\infty}^{\infty} d t f(t) F_{t} \mu \tag{4.48}
\end{equation*}
$$

is smooth for all $f \in L_{+}^{1}(\mathbb{R}, d t)$.

The finite weight $\psi$ is integrable if and only if $\mu_{\psi}$ is smooth. The smoothable measures are precisely the measures of the form $\int \mu_{b} d \nu_{B}(b)$ with respect to some base and ceiling function construction of the flow. It is then obvious from eqs. (4.45)-(4.47) that given any smoothable measure $\mu$, one can choose $h_{b}, b \in B$ so that $\phi(h f(h))=\mu(f), f \in L^{\infty}(X, \nu)$. I.e. every smoothable measure $\mu$ is of the form $\mu_{\psi}$ for some $\psi \in M_{*}^{+}$.

## 5. Comparison of finite periodic weights

Let $M$ be a type $\mathrm{III}_{0}$ factor with $T \in T(M)$ for some $0<T<\infty$, and let $\phi$ be a faithful state on $M, \sigma_{T}^{\phi}=1$. Let $\lambda, \tilde{M}, \tilde{\phi}, \tilde{S}, \tilde{M}_{\tilde{\phi}}, C_{\tilde{\phi}}, \theta, B, \nu_{B}$ be as in eqs. (4.1)-(4.12). To each finite weight $\psi \in \tilde{M}_{\boldsymbol{\phi}}^{T}$ (see definition 5.1) we associate a finite measure $\mu_{\psi}$ on $B, \mu_{\psi}<\nu_{B}$. We establish a number of properties of the map $\psi \rightarrow \mu_{\psi}$. In particular lemmas $5.7,5.8$ and 5.9 are required for the proof that AT implies ITPFI in the discrete case $T(M) \neq\{0\}$ (see lemma 8.2).

If $\psi$ is a weight on $M$ such that $\sigma_{T}^{\psi}=1$, then $(D \psi: D \phi)_{T}=e^{i \alpha} s(\psi)$ where $0 \leq \alpha<$ $2 \pi$, and ( $D \psi: D \phi)_{t}, t \in \mathbb{R}$ is the cocyle Radon-Nikodym derivative (see [5, pp. 478479]). The weight $\theta=\psi \oplus \phi$ will satisfy $\sigma_{T}^{\theta}=1$ if and only if $\alpha=0$. Note that for $\beta>0$ we have

$$
\begin{equation*}
(D(\beta \psi): D \phi)_{T}=\beta^{i T}(D \psi: D \phi)_{T} \tag{5.1}
\end{equation*}
$$

so that for some $\lambda<\beta \leq 1$ we have ( $D(\beta \psi): D \phi)_{T}=s(\psi)$.
Definition 5.1. Let $\omega$ be a faithful weight on the von Neumann algebra $\mathscr{A}$ such that $\sigma_{T}^{\omega}=1$ for some $0<T<\infty$. Then $\mathscr{A}_{\omega}^{T}$ denotes the set of all finite weights $\psi$ on $\mathscr{A}$ such that $(D \psi: D \omega)_{T}=s(\psi)$. If $u$ is a partial isometry in $\mathscr{A}$ with $u u^{*} \in \mathscr{A}_{\omega}$ then the equation $\psi(x)=\omega\left(u x u^{*}\right), x \in \mathscr{A}_{+}$defines a weight $\psi$ with support $u^{*} u$. We write $\psi=\omega_{u}$.
Definition 5.2. Let $\omega$ and $\psi$ be weights on a von Neumann algebra $\mathscr{A}$. We say that $\omega$ and $\psi$ are equivalent and write $\omega \sim \psi$ if there exists a partial isometry $u \in \mathscr{A}$ such that $u u^{*}=s(\omega), u^{*} u=s(\psi)$ and $\psi=\omega_{u}$.

Lemma 5.3. Let $M, T, \phi, \tilde{M}, \tilde{\phi}$ be as above. Let $\psi \in \tilde{M}_{\dot{\phi}}^{T}$. Then there exists a projection $e \in \tilde{M}_{\tilde{\phi}}$ such that $\psi \sim \tilde{\phi}_{e}$. If $e, f$ are projections in $\tilde{M}_{\tilde{\phi}}$ then $\tilde{\phi}_{e} \sim \tilde{\phi}_{f}$ if and only if $e \sim f$ in $\tilde{M}_{\dot{\phi}}$.
Proof. The assertion $\tilde{\phi}_{e} \sim \tilde{\phi}_{f}$ if and only if $e \sim f$ in $\tilde{M}_{\tilde{\phi}}$ is lemma 1.4(d) of [5]. To prove the first assertion, consider the weight $\theta$ defined on $P=\tilde{M} \otimes F_{2}$ by

$$
\begin{equation*}
\theta\left(\sum_{i, j=1}^{2} x_{i j} \otimes e_{i j}\right)=\tilde{\phi}\left(x_{11}\right)+\psi\left(x_{22}\right) . \tag{5.2}
\end{equation*}
$$

From ([2, lemma 1.2.2]) and ([5, pp.478-9]) we have

$$
\begin{align*}
\sigma_{t}^{\theta}\left(x \otimes e_{11}\right) & =\sigma_{t}^{\tilde{\delta}}(x) \otimes e_{11},  \tag{5.3}\\
\sigma_{t}^{\theta}\left(s(\psi) \otimes e_{21}\right) & =u_{t} \otimes e_{21}, \tag{5.4}
\end{align*}
$$

and

$$
\begin{equation*}
\sigma_{1}^{\theta}\left(x \otimes e_{22}\right)=\sigma_{1}^{\psi}(x) \otimes e_{22} \tag{5.5}
\end{equation*}
$$

for all $x \in M$, where $u_{t}=(D \psi: D \tilde{\phi})_{t}, t \in \mathbb{R}$. By [5, lemma 1.4(b)]) we have $\psi \sim \tilde{\phi}_{e}$
for some $e$ if and only if

$$
\begin{equation*}
s(\psi) \otimes e_{22}<1 \otimes e_{11}\left(P_{\theta}\right) \tag{5.6}
\end{equation*}
$$

(i.e. $s(\psi) \otimes e_{22}$ is equivalent in the centralizer $P_{\theta}$ to a sub-projection of $1 \otimes e_{11}$ ). Since $\tilde{M}_{\tilde{\phi}}$ is properly infinite (see eq. (4.10)), there exist projections $e_{j} \in \tilde{M}_{\tilde{\phi}}, j \in \mathbb{Z}$ such that $e_{j} e_{k}=0$ if $j \neq k, e_{j} \sim e_{k}$ and $\sum e_{j}=1$. It follows from eq. (5.3) that $e_{j} \otimes e_{11} \in P_{\theta}$, $e_{j} \otimes e_{11} \sim e_{k} \otimes e_{11}\left(P_{\theta}\right)$, and hence that $1 \otimes e_{11}=\sum_{j \in Z} e_{j} \otimes e_{11}$ is a properly infinite projection in $P_{\theta}$. Hence to prove eq. (5.6) it suffices to show that given any projection $f \leq s(\psi), f \neq 0, f \in \tilde{M}_{\psi}$ there exists $y \in P_{\theta}$ such that

$$
\begin{equation*}
\left(f \otimes e_{22}\right) y\left(1 \otimes e_{11}\right) \neq 0 \tag{5.7}
\end{equation*}
$$

Since $\sigma_{T}^{\theta}=1$ it follows from eq. (5.4) that $u_{t}$ is periodic with period $T$ and hence

$$
\begin{equation*}
u_{t}=\sum_{k \in \mathbb{Z}} u^{(k)} e^{i 2 \pi k t / T} \tag{5.8}
\end{equation*}
$$

where

$$
\begin{equation*}
u^{(k)} \otimes e_{21}=\frac{1}{T} \int_{0}^{T} d t \sigma_{t}^{\theta}\left(s(\psi) \otimes e_{21}\right) e^{-i 2 \pi k t / T} \tag{5.9}
\end{equation*}
$$

Now

$$
\begin{align*}
0 \neq \sigma_{t}^{\theta}\left(f \otimes e_{21}\right) & =\sigma_{t}^{\theta}\left(\left(f \otimes e_{22}\right)\left(s(\psi) \otimes e_{21}\right)\right) \\
& =f u_{t} \otimes e_{21} \tag{5.10}
\end{align*}
$$

and hence $f u^{(K)} \neq 0$ for some $K \in \mathbb{Z}$. Eq. (5.9) gives

$$
\begin{equation*}
\sigma_{t}^{\theta}\left(u^{(K)} \otimes e_{21}\right)=e^{i 2 \pi K t / T} u^{(K)} \otimes e_{21} \tag{5.11}
\end{equation*}
$$

It follows from eq. (4.3) that the unitary $\tilde{S}$ defined by eq. (4.8) satisfies

$$
\begin{equation*}
\sigma_{t}^{\tilde{\theta}}(\tilde{S})=e^{-i 2 \pi t / T} \tilde{S} \tag{5.12}
\end{equation*}
$$

since $\lambda^{i t}=e^{-i 2 \pi t / T}$. Define

$$
\begin{equation*}
y=\left(u^{(K)} \otimes e_{21}\right)\left(\tilde{S}^{-K} \otimes e_{11}\right) \tag{5.13}
\end{equation*}
$$

Eqs. (5.11)-(5.13) give $\sigma_{t}^{\boldsymbol{\theta}}(y)=y$, hence $y \in P_{\theta}$. Since $\tilde{S}$ is unitary and $f u^{(K)} \neq 0$, eq. (5.7) is satisfied.

Definition 5.4. Let $\tilde{M}, \tilde{\phi}$ be as above, $\psi \in \tilde{M}_{\phi}^{T}$. We define the measure $\mu_{\psi}$ associated with $\psi$ as the measure $\mu_{e}$ defined by eq. (4.15) where $\psi \sim \tilde{\phi}_{e}$.
Lemma 5.5. Let $\psi_{1}, \psi_{2} \in \tilde{M}_{\dot{\phi}}^{\tau}$. Then $\mu_{\psi_{1}}=\mu_{\psi_{2}}$ if and only if $\psi_{1} \sim \psi_{2}$.
Proof. The lemma follows immediately from lemma 5.3 and the fact that $\mu_{e}=\mu_{f}$ if and only if $e \sim f$ in $\tilde{M}_{\tilde{\Phi}}$ (see eq. (4.16)).
Lemma 5.6. Let $\psi_{1}, \psi_{2} \in \tilde{M}_{\phi}^{T}$ be such that $s\left(\psi_{1}\right) s\left(\psi_{2}\right)=0$. Then $\mu_{\psi_{1}+\psi_{2}}=\mu_{\psi_{1}}+\mu_{\psi_{2}}$.
Proof. Note that eq. (4.16) implies that if $e=e_{1}+e_{2}$ where $e, e_{1}, e_{2} \in \tilde{M}_{\dot{\phi}}^{+}$, then $\mu_{e}=\mu_{e_{1}}+\mu_{e_{2}}$. By lemma 5.3 we have $\psi_{j}=\tilde{\phi}_{u_{j}}, j=1,2$ where, since $\tilde{M}_{\tilde{\phi}}$ is properly infinite, we can choose $u_{1}$ and $u_{2}$ such that $e_{1} e_{2}=0$ where $e_{j}=u_{j}^{*} u_{j}$. Then $u=u_{1}+u_{2}$ is a partial isometry such that $\psi_{1}+\psi_{2}=\tilde{\phi}_{u}$, and $u^{*} u=e_{1}+e_{2}$. Since $\tilde{\phi}_{u} \sim \tilde{\phi}_{u^{*} u}$ and $\tilde{\phi}_{u_{j}} \sim \tilde{\phi}_{e_{j}}$ the result follows.
Lemma 5.7. Let $\psi \in \tilde{M}_{\phi}^{T}$. Let $\mu_{1}, \ldots, \mu_{n}$ be measures on $B$ such that $\mu_{\psi}=$ $\sum_{j=1}^{n} \mu_{j}$. Then there exist orthogonal projections $e_{1}, \ldots, e_{n}$ such that $s(\psi)=\sum_{j=1}^{n} e_{j}$ and $\mu_{\psi_{e_{j}}}=\mu_{j}$.

Proof. Since every $\mu<\nu_{B}$ occurs as $\mu_{\rho}$ for some $\rho \in \tilde{M}_{\phi}^{T}$ (see eq. (4.33)), there exist mutually orthogonal projections $f_{1}, \ldots, f_{n} \in \tilde{M}$ and $\rho_{1}, \ldots, \rho_{n} \in \tilde{M}_{\phi}^{T}$ such that

$$
\begin{equation*}
s\left(\rho_{j}\right)=f_{j} \tag{5.14}
\end{equation*}
$$

and

$$
\begin{equation*}
\mu_{\rho_{j}}=\mu_{j}, \quad j=1, \ldots, n . \tag{5.15}
\end{equation*}
$$

By lemma 5.5 and $5.6, \rho=\sum \rho_{j} \sim \psi$. Hence there is a partial isometry $u$ with $u^{*} u=s(\rho)=\sum f_{j}, u u^{*}=s(\psi)$, and $\psi=\rho_{u}$. Then $e_{j}=u f_{j}$ are the desired projections.
Lemma 5.8. Let $\psi \in \tilde{M}_{\boldsymbol{\phi}}^{T}, \psi \sim \tilde{\phi}_{e}$. Then

$$
\theta \mu_{\psi}=\lambda \mu_{\lambda^{-1} \psi}=\lambda \mu_{\theta(e)} .
$$

Proof. We have $\mu_{\psi}=\mu_{e}$. Let $c \in C_{\bar{\phi}}$. Using eqs. (4.9) and (4.15) we have

$$
\begin{align*}
\left(\theta \mu_{\psi}\right)(c) & =\mu_{e}\left(\theta^{-1}(c)\right)=\tilde{\phi}\left(e \theta^{-1}(c)\right) \\
& =(\tilde{\phi} \circ \theta)(\theta(e) c)=\lambda \tilde{\phi}(\theta(e) c)=\lambda \mu_{\theta(e)}(c) \tag{5.16}
\end{align*}
$$

It remains only to prove that $\mu_{\lambda \psi}=\mu_{\theta^{-1}(e)}$. Let $u$ be a partial isometry in $\tilde{M}$ such that $\psi=\tilde{\phi}_{u}$ and $u u^{*}=e$. Then

$$
\begin{equation*}
(\lambda \psi)(x)=\lambda \tilde{\phi}\left(u x u^{*}\right)=\tilde{\phi}\left(\theta^{-1}\left(u x u^{*}\right)\right), \quad x \in \tilde{M} \tag{5.17}
\end{equation*}
$$

Since $\theta=\operatorname{Ad} \tilde{S}$ we get

$$
\begin{equation*}
(\lambda \psi)(x)=\tilde{\phi}\left(\tilde{S} u x u^{*} \tilde{S}^{*}\right)=\tilde{\phi}_{\tilde{S}_{u}}(x) . \tag{5.18}
\end{equation*}
$$

But $\tilde{\phi}_{\omega} \sim \tilde{\phi}_{\omega \omega^{*}}$ and $\tilde{S} u u^{*} \tilde{S}^{*}=\theta^{-1}(e)$.
Lemma 5.9. Let $\psi_{1} \in \tilde{M}_{\Phi}^{\tau}, \psi_{1} \neq 0$, and $\varepsilon>0$. Let $\mu_{2} \neq 0$ be a finite measure on $B$, $\mu_{2}<\nu_{B}$. Then there exists $\psi_{2} \in \tilde{M}_{\Phi}^{T}$ such that
(i) $s\left(\psi_{2}\right)=s\left(\psi_{1}\right)$;
(ii) $\mu_{\psi_{2}}=\mu_{2}$; and
(iii) $\left\|\psi_{1}-\psi_{2}\right\| \leq\left\|\mu_{\psi_{1}}-\mu_{2}\right\|+\varepsilon$.

Froof. Write $\mu_{1}=\mu_{\psi_{1}}$ and $f_{j}=d \mu_{j} / d \nu_{B}, j=1,2$. By a routine argument one can choose a family of projections $e(\alpha) \in R_{0,1}, 0<\alpha<\infty$ such that

$$
\begin{equation*}
e(\alpha) e(\beta)=e(\alpha) \quad \text { if } \alpha<\beta \tag{5.19}
\end{equation*}
$$

and

$$
\begin{equation*}
\tau(e(\alpha))=\alpha \tag{5.20}
\end{equation*}
$$

where $\tau$ is the trace on $R_{0,1}$ such that $\tilde{\phi}=\tau \otimes \nu_{\mathrm{B}}$ (see eqs. (4.11) and (4.12)). Then $\mu_{j}=\mu_{e_{j}} j=1,2$ where

$$
\begin{equation*}
e_{j}=\int_{B}^{\oplus} e\left(f_{j}(b)\right) d \nu_{B}(b) \tag{5.21}
\end{equation*}
$$

Let

$$
\begin{align*}
B_{+} & =\left\{b: f_{2}(b)>f_{1}(b)\right\}  \tag{5.22}\\
B_{-} & =\left\{b: f_{2}(b)<f_{1}(b)\right\}  \tag{5.23}\\
B_{0} & =\left\{b: f_{2}(b)=f_{1}(b)\right\} \tag{5.24}
\end{align*}
$$

Then

$$
\begin{equation*}
e_{j}=e_{j+}+e_{j-}+e_{j 0} \tag{5.25}
\end{equation*}
$$

where

$$
\begin{equation*}
e_{j+}=\int_{B_{+}}^{\oplus} e\left(f_{j}(b)\right) d \nu_{B}(b) \tag{5.26}
\end{equation*}
$$

etc. Let $u$ be a partial isometry in $\tilde{M}$ such that $\psi_{1}=\tilde{\phi}_{u}$ where $u u^{*}=e_{1}, u^{*} u=s\left(\psi_{1}\right)$. Case (i). $\nu_{B}\left(B_{+}\right)=\nu_{B}\left(B_{-}\right)=0$. Take $\psi_{2}=\psi_{1}$.
Case (ii). $\nu_{B}\left(B_{+}\right), \nu_{B}\left(B_{-}\right)>0$. Let $\omega$ be a partial isometry in $\tilde{M}$ mapping the non-zero projection $e_{2+}-e_{1+}$ onto the non-zero projection $e_{1-}-e_{2}$. Define $\psi_{2}=\psi_{u}$ where

$$
\begin{equation*}
\psi=\tilde{\phi}_{e_{1}-\left(e_{1-}-e_{2-}\right)}+\left(\tilde{\phi}_{e_{2+}-e_{1+}}\right)_{\omega} . \tag{5.27}
\end{equation*}
$$

Then $\psi \sim \dot{\phi}_{e_{2}}$ so that $\mu_{\psi}=\mu_{2}$. Furthermore $s(\psi)=e_{1}, s\left(\psi_{2}\right)=u^{*} u=s\left(\psi_{1}\right)$, and

$$
\begin{equation*}
\left\|\psi_{2}-\psi_{1}\right\|=\left\|\psi-\tilde{\phi}_{e_{1}}\right\| . \tag{5.28}
\end{equation*}
$$

We have

$$
\begin{equation*}
\psi-\tilde{\phi}_{e_{1}}=\left(\tilde{\phi}_{e_{2+}-e_{1+}}\right)_{\omega}-\tilde{\phi}_{e_{1-}-e_{2}} . \tag{5.29}
\end{equation*}
$$

Since $\left\|\tilde{\phi}_{f}\right\|=\tilde{\phi}(f)=\mu_{f}(B)$ we obtain

$$
\begin{align*}
\left\|\psi_{2}-\psi_{1}\right\| & \leq\left\|\tilde{\phi}_{e_{2+}-e_{1+}}\right\|+\left\|\tilde{\phi}_{e_{1-}-e_{2-}}\right\| \\
& =\left(\mu_{2}-\mu_{1}\right)\left(B_{+}\right)+\left(\mu_{1}-\mu_{2}\right)\left(B_{-}\right)=\left\|\mu_{1}-\mu_{2}\right\| . \tag{5.30}
\end{align*}
$$

Case (iii). $\nu_{B}\left(B_{+}\right)=0, \nu_{B}\left(B_{-}\right) \neq 0$. Choose $g(b), b \in B$ such that

$$
\begin{equation*}
0 \leq g(b) \leq f_{2}(b) \tag{5.31}
\end{equation*}
$$

and

$$
\begin{equation*}
0<\int_{B} g(b) d \nu_{B}(b)<\frac{1}{2} \varepsilon . \tag{5.32}
\end{equation*}
$$

Let

$$
\begin{equation*}
g=\int_{B}^{\oplus} e(g(b)) d \nu_{B}(b) . \tag{5.33}
\end{equation*}
$$

Let $\omega$ be a partial isometry in $\tilde{M}$ such that $\omega^{*} \omega=g, \omega \omega^{*}=e_{1}-e_{2}+g$. Define $\psi_{2}=\psi_{u}$ where

$$
\begin{equation*}
\psi=\tilde{\phi}_{e_{2}-g}+\left(\tilde{\phi}_{g}\right) \omega \tag{5.34}
\end{equation*}
$$

Then $\psi \sim \tilde{\phi}_{e_{2}}$ so that $\mu_{\psi}=\mu_{2}$. Furthermore $s(\psi)=e_{1}$ so that $s\left(\psi_{2}\right)=u^{*} u=s\left(\psi_{1}\right)$, and

$$
\begin{equation*}
\left\|\psi_{2}-\psi_{1}\right\|=\left\|\psi-\tilde{\phi}_{e_{1}}\right\| . \tag{5.35}
\end{equation*}
$$

Since

$$
\begin{equation*}
\tilde{\phi}_{e_{1}}-\psi=\tilde{\phi}_{g}+\tilde{\phi}_{e_{1}-e_{2}}-\left(\tilde{\phi}_{g}\right)_{\omega} \tag{5.36}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\|\tilde{\phi}_{g}\right\|=\tilde{\phi}(g)=\int g(b) d \nu_{B}<\frac{1}{2} \varepsilon \tag{5.37}
\end{equation*}
$$

we obtain $\left\|\psi_{2}-\psi_{1}\right\| \leq\left\|\mu_{1}-\mu_{2}\right\|+\varepsilon$.
Case (iv). $\nu_{B}\left(B_{+}\right) \neq 0, \nu_{B}\left(B_{-}\right)=0$. The argument is similar to case (iii).
6. Comparison of finite weights: the general case

We extend the results of $\S 5$ to the general case. This is straightforward except for lemma 5.9 , where we now use the lacunary construction (see lemma 6.4). In this section $M$ is a type $\mathrm{III}_{0}$ factor, and we follow the notation of $\S 4$.

Lemma 6.1. Let $\psi \in M_{*}^{+}$, then $\mu_{\lambda \psi}=\lambda \mathscr{F}_{\lambda^{-1}}^{M} \mu_{\psi}$. If $\psi_{1}, \psi_{2} \in M_{*}^{+}$then $\psi_{1} \sim \psi_{2} \Leftrightarrow \mu_{\psi_{1}}=\mu_{\psi_{2}}$, and $\mu_{\psi_{1}+\psi_{2}}=\mu_{\psi_{1}}+\mu_{\psi_{2}}$ if $s\left(\psi_{1}\right) s\left(\psi_{2}\right)=0$.
Proof: This is corollary 1.13 (ii) of [5].
Lemma 6.2. Let $\psi \in M_{*}^{+}, \mu_{\psi}=\sum_{j=1}^{n} \mu_{j}$. Then there exist orthogonal projections $e_{1}, \ldots, e_{n}$ such that $s(\psi)=\sum_{j=1}^{n} e_{j}$ and $\psi=\sum_{j=1}^{n} \psi_{e_{j}}$ where $\mu_{\psi_{e_{j}}}=\mu_{j}$.
Proof. Since every smoothable measure occurs as $\mu_{\chi}$ for some $\chi$ (see eq. (4.48) et seq.) the proof of lemma 5.7 holds verbatim.

The next lemma is a technical result needed in the proof of lemma 6.4.
Lemma 6.3. Let $\sigma_{1}, \sigma_{2}$ be non-atomic measures on $I=[0,1], \sigma_{1}(I)=\sigma_{2}(I)<\infty$. Let $F_{j}(x)=\sigma_{j}([0, x])$, and let

$$
S_{j}=\left\{x \in I: F_{j}(y)=F_{j}(x) \text { implies } y=x\right\}
$$

$j=1$, 2. Then $\sigma_{j}\left(I \backslash S_{j}\right)=0$ and the equation $F_{1}(\gamma(x))=F_{2}(x)$ defines a monotonic bijection $\gamma: S_{1} \rightarrow S_{2}$ satisfying

$$
\int_{0}^{1} f(\gamma(t)) d \sigma_{1}(t)=\int_{0}^{1} f(t) d \sigma_{2}(t)
$$

for all $f \in L^{1}\left(I, \sigma_{2}\right)$. Furthermore

$$
\begin{equation*}
\left\|\sigma_{1}-\sigma_{2}\right\| \geq \int_{0}^{1}|t-\gamma(t)| d \sigma_{1}(t) \tag{6.1}
\end{equation*}
$$

Proof. To prove eq. (6.1) consider the function $f$ defined by

$$
f^{\prime}(t)=\left\{\begin{align*}
+1 & \text { if } t>\gamma(t)  \tag{6.2}\\
0 & \text { if } t=\gamma(t) \\
-1 & \text { if } t<\gamma(t)
\end{align*}\right.
$$

and $f(0)=0$. Since $\|f\|_{\infty} \leq 1$ we have

$$
\begin{align*}
\left\|\sigma_{1}-\sigma_{2}\right\| & \geq\left|\sigma_{1}(f)-\sigma_{2}(f)\right| \\
& =\left|\int[f(t)-f(\gamma(t))] d \sigma_{1}(t)\right| \\
& =\int|t-\gamma(t)| d \sigma_{1}(t) \tag{6.3}
\end{align*}
$$

All other properties are obvious.
Lemma 6.4. Let $\psi_{1}$ be a finite integrable weight on $M$, and let $\mu_{2}$ be a smooth measure on $X$ (i.e. $\mu_{2}<\nu$, see definition 4.2). Then there exists $\psi_{2} \in M_{*}^{+}$such that:
(i) $s\left(\psi_{2}\right)=s\left(\psi_{1}\right)$;
(ii) $\mu_{\psi_{2}}=\mu_{2} ;$ and
(iii) $\left\|\psi_{1}-\psi_{2}\right\| \leq 5\left\|\mu_{\psi_{1}}-\mu_{2}\right\|$.

Proof. We use the lacunary state construction of the flow of weights (eqs. (4.39)(4.48)) where we choose the lacunary state $\phi$ so that $\frac{1}{2} \leq \rho<1$. We have

$$
\begin{equation*}
\psi_{1}(x)=\phi\left(h_{1} E\left(u x u^{*}\right)\right), \tag{6.4}
\end{equation*}
$$

where $h_{1} \in M_{\phi}^{+}$so that

$$
\begin{equation*}
h_{1}=\int_{B}^{\oplus} d \nu_{B}(b) \int_{\rho(b)}^{1_{\oplus}} t d \sigma_{1, b}(t) \tag{6.5}
\end{equation*}
$$

where we can choose the measures $\sigma_{1, b}$ so that for $f \in L^{1}\left(X, \mu_{1}\right)$ we have

$$
\begin{equation*}
\mu_{1}(f)=\int_{B} d \nu_{B}(b) \int_{\rho(b)}^{1} t f(b, t) d \sigma_{1, b}(t) \tag{6.6}
\end{equation*}
$$

where $\mu_{1}=\mu_{\psi_{1}}$. By the remark following definition 4.2, we can choose measures $\sigma_{2, b}$ such that

$$
\begin{equation*}
\mu_{2}(f)=\int_{B} d \nu_{B}(b) \int_{\rho(b)}^{1} t f(b, t) d \sigma_{2, b}(t) \tag{6.7}
\end{equation*}
$$

We begin by altering the measures slightly so that lemma 6.3 can be used. Let

$$
\begin{equation*}
B_{+}=\left\{b \in B: \sigma_{1, b}([\rho(b), 1]) \geq \sigma_{2, b}([\rho(b), 1])\right\} \tag{6.8}
\end{equation*}
$$

and

$$
\begin{equation*}
B_{-}=B \backslash B_{+} \tag{6.9}
\end{equation*}
$$

For $b \in B_{+}$define $t_{2, b}=1$ and

$$
\begin{equation*}
t_{1, b}=\sup \left\{\rho(b) \leq t \leq 1: \sigma_{1, b}([\rho(b), t])=\sigma_{2, b}([\rho(b), 1])\right\} \tag{6.10}
\end{equation*}
$$

and for $b \in B_{-}$define $t_{1, b}=1$ and

$$
\begin{equation*}
t_{2, b}=\sup \left\{\rho(b) \leq t \leq 1: \sigma_{2, b}([\rho(b), t])=\sigma_{1, b}([\rho(b), 1])\right\} \tag{6.11}
\end{equation*}
$$

Define $\mu_{1}^{\prime}, \mu_{2}^{\prime}$ by

$$
\begin{equation*}
\left(d \mu_{j}^{\prime} / d \mu_{j}\right)(b, t)=\chi_{\left[\rho(b), t_{i, f}\right]}(t) \tag{6.12}
\end{equation*}
$$

Then

$$
\begin{align*}
\left\|\mu_{1}-\mu_{1}^{\prime}\right\| & =\mu_{1}\left(\chi_{\left[t_{1,5,1]}\right]}\right) \\
& \leq \nu_{B} \circ \sigma_{1}\left(\chi_{\left[t_{1}, b_{r} 1\right]}\right)=\left(\nu_{B} \circ \sigma_{1}-\nu_{B} \circ \sigma_{2}\right)\left(\chi_{B_{+}}\right) \\
& \leq\left\|\nu_{B} \circ \sigma_{1}-\nu_{B} \circ \sigma_{2}\right\| \leq 2\left\|\mu_{1}-\mu_{2}\right\|, \tag{6.13}
\end{align*}
$$

where the last inequality follows from the fact that $d \mu_{j} / d \nu_{B} \circ \sigma_{j}=t \geq \frac{1}{2}$. Similarly

$$
\begin{equation*}
\left\|\mu_{2}-\mu_{2}^{\prime}\right\| \leq 2\left\|\mu_{1}-\mu_{2}\right\| . \tag{6.14}
\end{equation*}
$$

We can now write $\mu_{j}^{\prime}=\nu_{B} \circ \sigma_{j}^{\prime}, j=1,2$ where

$$
\begin{equation*}
\sigma_{1, b}^{\prime}([\rho(b), 1])=\sigma_{2, b}^{\prime}([\rho(b), 1]), \quad b \in B \tag{6.15}
\end{equation*}
$$

By lemma 6.3 there exists a measurable function $\gamma(b, t)$ such that for all $f \in L^{\prime}\left(X, \mu_{1}^{\prime}\right)$ we have

$$
\begin{equation*}
\int_{\rho(b)}^{1} \gamma(b, t) f(b, \gamma(b, t)) d \sigma_{1, b}^{\prime}(t)=\int_{\rho(b)}^{1} t f(b, t) d \sigma_{2, b}^{\prime}(t) \tag{6.16}
\end{equation*}
$$

for a.e. $b \in B$. We define

$$
\begin{equation*}
h_{1}^{\prime}=\int_{B}^{\oplus} d \nu_{B}(b) \int_{\rho(b)}^{1_{\oplus}} t d \sigma_{1, b}^{\prime}(t) \tag{6.17}
\end{equation*}
$$

$$
\begin{equation*}
h_{2}^{\prime}=\int_{B}^{\oplus} d \nu_{B}(b) \int_{\rho(b)}^{1_{\oplus}} \gamma(b, t) d \sigma_{1, b}^{\prime}(t), \tag{6.18}
\end{equation*}
$$

and

$$
\begin{equation*}
\psi_{j}^{\prime}(x)=\phi\left(h_{j}^{\prime} E\left(u x u^{*}\right)\right), \quad j=1,2 . \tag{6.19}
\end{equation*}
$$

Using eq. (6.13) we obtain

$$
\begin{equation*}
\left\|\psi_{1}-\psi_{1}^{\prime}\right\|=\phi\left(\left|h_{1}-h_{1}^{\prime}\right|\right)=\left\|\mu_{1}-\mu_{1}^{\prime}\right\| \leq 2\left\|\mu_{1}-\mu_{2}\right\| . \tag{6.20}
\end{equation*}
$$

We also have

$$
\begin{align*}
\left\|\psi_{1}^{\prime}-\psi_{2}^{\prime}\right\| & =\phi\left(\left|h_{1}^{\prime}-h_{2}^{\prime}\right|\right) \\
& =\int d \nu_{B}(b) \int|t-\gamma(b, t)| d \sigma_{1, b}^{\prime}(t) \\
& \leq \int d \nu_{B}(b)\left\|\sigma_{1, b}^{\prime}-\sigma_{2, b}^{\prime}\right\|=\left\|\mu_{1}^{\prime}-\mu_{2}^{\prime}\right\| \tag{6.21}
\end{align*}
$$

where the last inequality follows from lemma 6.3.
Case (i). $\mu_{1}^{\prime}=\mu_{1}$ and $\mu_{2}^{\prime}=\mu_{2}$. Then $\psi_{2}=\psi_{2}^{\prime}$ satisfies the lemma.
Case (ii). $\mu_{1} \neq \mu_{1}^{\prime}$ and $\mu_{2} \neq \mu_{2}^{\prime}$. Let $h_{1}^{\prime \prime}=h_{1}-h_{1}^{\prime}$. Then $s\left(h_{1}^{\prime \prime}\right) s\left(h_{1}^{\prime}\right)=0$ and $s\left(h_{1}^{\prime \prime}\right) \neq 0$.
Construct $h_{2}^{\prime \prime}$ such that $s\left(h_{2}^{\prime \prime}\right)=s\left(h_{1}^{\prime \prime}\right)$ and $\mu_{h_{2}^{\prime \prime}}=\mu_{2}-\mu_{2}^{\prime}$. Let $h_{2}=h_{2}^{\prime}+h_{2}^{\prime \prime}$ and define

$$
\begin{equation*}
\psi_{2}=\phi\left(h_{2} E\left(u x u^{*}\right)\right) \tag{6.22}
\end{equation*}
$$

By construction $\mu_{\psi_{2}}=\mu_{2}$ and $s\left(\psi_{2}\right)=s\left(\psi_{1}\right)$. We have $\psi_{2}=\psi_{2}^{\prime}+\psi_{2}^{\prime \prime}$ where

$$
\begin{equation*}
\left\|\psi_{2}^{\prime \prime}\right\|=\phi\left(\left|h_{2}^{\prime \prime}\right|\right)=\left\|\mu_{2}-\mu_{2}^{\prime}\right\| \leq 2\left\|\mu_{1}-\mu_{2}\right\| \tag{6.23}
\end{equation*}
$$

(see eq. (6.14)). Eqs. (6.20), (6.21) and (6.23) give

$$
\begin{equation*}
\left\|\psi_{1}-\psi_{2}\right\| \leq 5\left\|\mu_{1}-\mu_{2}\right\| . \tag{6.24}
\end{equation*}
$$

Case (iii). $\mu_{1}=\mu_{1}^{\prime}$ and $\mu_{2} \neq \mu_{2}^{\prime}$. Choose a projection $e \in M$ such that $0<e<s\left(h_{2}^{\prime}\right)$, and

$$
\begin{equation*}
\phi\left(e h_{2}^{\prime}\right) \leq\left\|\mu_{1}-\mu_{2}\right\| . \tag{6.25}
\end{equation*}
$$

Choose $h_{2}^{\prime \prime}$ such that $s\left(h_{2}^{\prime \prime}\right)=e$ and

$$
\begin{equation*}
\mu_{h_{2}^{\prime \prime}}=\mu_{2}-\mu_{2}^{\prime}+\mu_{e h_{2}^{\prime}} . \tag{6.26}
\end{equation*}
$$

Then eqs. (6.14) and (6.25) give

$$
\begin{equation*}
\left\|\mu_{h_{2}^{2}}\right\| \leq 3\left\|\mu_{1}-\mu_{2}\right\| \tag{6.27}
\end{equation*}
$$

Define

$$
\begin{equation*}
\psi_{2}(x)=\phi\left(\left(h_{2}^{\prime \prime}+(1-e) h_{2}^{\prime}\right) E\left(u x u^{*}\right)\right) \tag{6.28}
\end{equation*}
$$

By construction $\mu_{\psi_{2}}=\mu_{2}$ and $s\left(\psi_{2}\right)=s\left(\psi_{1}\right)$. Eqs. (6.25) and (6.27) give

$$
\begin{equation*}
\left\|\psi_{2}-\psi_{2}^{\prime}\right\| \leq 4\left\|\mu_{1}-\mu_{2}\right\| . \tag{6.29}
\end{equation*}
$$

Since $\psi_{1}=\psi_{1}^{\prime}$, eq. (6.21) now gives

$$
\begin{equation*}
\left\|\psi_{1}-\psi_{2}\right\| \leq 5\left\|\mu_{1}-\mu_{2}\right\| . \tag{6.30}
\end{equation*}
$$

Case (iv). $\mu_{1} \neq \mu_{1}^{\prime}$ and $\mu_{2}=\mu_{2}^{\prime}$. Choose a projection $e \in M$ such that $0<e<s\left(h_{1}^{\prime}\right)=$ $s\left(h_{2}^{\prime}\right)$ and

$$
\begin{equation*}
\phi\left(e h_{2}^{\prime}\right) \leq\left\|\mu_{1}-\mu_{2}\right\| . \tag{6.31}
\end{equation*}
$$

Choose $h_{2}^{\prime \prime}$ such that $s\left(h_{2}^{\prime \prime}\right)=s\left(h_{1}\right)-s\left(h_{1}^{\prime}\right)+e$ and

$$
\begin{equation*}
\mu_{h_{2}^{n}}=\mu_{e h_{2}^{\prime}} . \tag{6.32}
\end{equation*}
$$

Define $\psi_{2}$ by eq. (6.28). By construction $\mu_{\psi_{2}}=\mu_{2}$ and $s\left(\psi_{2}\right)=s\left(\psi_{1}\right)$. Eqs. (6.31), (6.32) and (6.28) give

$$
\begin{equation*}
\left\|\psi_{2}-\psi_{2}^{\prime}\right\| \leq 2\left\|\mu_{1}-\mu_{2}\right\| . \tag{6.33}
\end{equation*}
$$

Eqs. (6.20) and (6.21) now give

$$
\begin{equation*}
\left\|\psi_{1}-\psi_{2}\right\| \leq 5\left\|\mu_{1}-\mu_{2}\right\| . \tag{6.34}
\end{equation*}
$$

## 7. Product property and ITPFI factors

We introduce the 'product property' (definition 7.1) which is a variation of Størmer's property of being 'asymptotically a product state' [19]. This is a technical property which is equivalent to ITPFI (corollary 7.4 and lemma 7.6). Its purpose is to simplify the task of verifying the ITPFI property by eliminating the iterative part of the argument.
Definition 7.1. Let $M$ be a von Neumann algebra. A finite weight $\phi$ on $M$ is said to have the product property if given $\varepsilon>0$, a strong neighbourhood $V$ of 0 , and $x_{1}, \ldots, x_{n} \in M$, there exists a finite type I factor $K \subset M$ and finite weights $\phi_{1}, \phi_{2}$ on $K, K^{c}=K^{\prime} \cap M$ respectively such that:
(i) $x_{j} \in K+V, j=1, \ldots, n$, and
(ii) $\left\|\phi-\phi_{1} \otimes \phi_{2}\right\|<\varepsilon$.

If $M$ has a faithful finite weight with the product property, then $M$ is said to have the product property.
Remark 7.2. It follows immediately from the above definition that the product property implies approximately type I (see [8]). Clearly one can require that $\phi_{1}$ and $\phi_{2}$ are faithful, and (ii) can be replaced by
(ii) $\left\|\phi-\left.\left.\phi\right|_{K} \otimes \phi\right|_{K^{c}}\right\|<\varepsilon$.

If one were to study von Neumann algebras of the form $\bigotimes_{\nu}\left(\mathscr{A}_{\nu}, \phi_{\nu}\right)$ where the $\mathscr{A}_{\nu}$ are finite-dimensional matrix algebras, the appropriate product property would be to require only that $K$ be a finite-dimensional subalgebra.

In order to prove that an ITPFI factor has the product property we will use the following martingale condition, which was introduced by Araki and Woods ([1, lemma 6.10]).
Lemma 7.3. Let $M=\bigotimes_{k=1}^{\infty}\left(M_{k}, \phi_{k}\right)$ be an ITPFI factor. For each $n \in \mathcal{N}$ there is a conditional expectation $E_{n}: M \rightarrow M^{(n)}=\otimes_{k=1}^{n} M_{k}$ such that
(i) $E_{n}(x) \rightarrow x$ strongly for all $x \in M$;
(ii) $E_{n}\left(M_{\phi}\right)=M_{\phi}^{(n)}(n)$ where $\phi=\bigotimes_{k=1}^{\infty} \phi_{k}$ and $\phi^{(n)}=\bigotimes_{k=1}^{n} \phi_{k}$; and
(iii) $E_{n} E_{m}=E_{n}$ if $n<m$.

Proof. Define $E_{n}$ by the equation

$$
\begin{equation*}
\left(E_{n}(x) \alpha, \beta\right)=\left(x\left(\alpha \otimes\left(\bigotimes_{k=n+1}^{\infty} \Phi_{k}\right)\right), \beta \otimes\left(\bigotimes_{k=n+1}^{\infty} \Phi_{k}\right)\right) \tag{7.1}
\end{equation*}
$$

for all $x \in M, \alpha, \beta \in \otimes_{k=1}^{n} H_{k}$ where $M_{k}$ acts on $H_{k}$ and $\phi_{k}(y)=\left(y \Phi_{k}, \Phi_{k}\right)$ for all $y \in M_{k}$. That $E_{n}$ is a conditional expectation and properties (i) and (iii) follow directly from eq. (7.1) (compare [1, lemma 6.10]). Condition (ii) follows from

$$
\begin{equation*}
\sigma_{t}^{\phi^{(n)}}\left(E_{n}(x)\right)=E_{n}\left(\sigma_{t}^{\phi}(x)\right) \tag{7.2}
\end{equation*}
$$

which in turn follows from the observation that

$$
\begin{equation*}
\Delta_{\phi}=\Delta_{\phi^{(n)}} \otimes \Delta_{\left(\otimes \otimes_{k} \underline{x+1} \phi_{k}\right)} \tag{7.3}
\end{equation*}
$$

and a routine calculation using eq. (7.1).
Corollary 7.4. Any ITPFI factor has the product property.
Proof. Lemma 7.3 implies that any product state $\otimes \phi_{k}$ has the product property.
Lemma 7.5. Let $M$ be a von Neumann algebra with the product property. Then any finite weight $\psi$ has the product property.
Proof. Let $\varepsilon>0, V$ a strong neighbourhood of 0 , and $x_{1}, \ldots, x_{n} \in M$. Let $\phi$ be a faithful finite weight on $M$ with the product property. It follows from the HahnBanach theorem that the set of all finite weights $\chi$ such that $\chi \leq \lambda \phi$ for some $\lambda>0$, is norm-dense in $M_{*}^{+}$. It then follows from [17, theorem $1.24 .3, \mathrm{p} .76$ ] that there exists $h \in M^{+}$such that

$$
\begin{equation*}
|\psi(a)-\phi(h a h)| \leq \varepsilon\|a\|, \tag{7.4}
\end{equation*}
$$

for all $a \in M$. By assumption there exist a finite type I factor $K \subset M, k \in K$, and finite weights $\phi_{1}, \phi_{2}$ such that

$$
\begin{gather*}
x_{j} \in K+V, \quad j=1, \ldots, n  \tag{7.5}\\
\phi\left(|h-k|^{2}\right) \leq \varepsilon^{2}\|h\|^{-2}\|\phi\|^{-1},  \tag{7.6}\\
\left\|\phi-\phi_{1} \otimes \phi_{2}\right\| \leq \varepsilon\|h\|^{-2}, \tag{7.7}
\end{gather*}
$$

and

$$
\begin{equation*}
\|k\| \leq\|h\| \tag{7.8}
\end{equation*}
$$

That one can achieve eq. (7.8) follows, as in the proof of the Kaplansky density theorem, by approximating an element $h^{\prime} \in M$ such that $h=2 h^{\prime}\left(1+\left(h^{\prime}\right)^{2}\right)^{-1}$ (see [7, p. 44]). Since

$$
\phi(h a h)=\phi(k a k)+\phi((h-k) a h)+\phi(k a(h-k)),
$$

it follows from the Cauchy-Schwarz inequality and eq. (7.6) that

$$
\begin{equation*}
|\phi(h a h)-\phi(k a k)| \leq 2 \varepsilon\|a\| . \tag{7.9}
\end{equation*}
$$

Eqs. (7.4), (7.7) and (7.9) give

$$
\begin{equation*}
\left\|\psi-\psi_{1} \otimes \psi_{2}\right\| \leq 4 \varepsilon, \tag{7.10}
\end{equation*}
$$

where $\psi_{1}(a)=\phi_{1}(k a k), a \in K$, and $\psi_{2}(a)=\phi_{2}(a), a \in K^{c}$.
Lemma 7.6. Let $M$ be a properly infinite von Neumann algebra with the product property. Then $M$ is ITPFI.
Proof. Let $\psi$ be a faithful state on $M$. Let $\left(x_{j}\right)_{j \in \mathcal{N}}$ be dense in the unit ball $M_{1}$ of $M$, let $V_{j}$ be a sequence of strong neighbourhoods of 0 decreasing to 0 , and let
$\eta_{n}>0$ with

$$
\begin{equation*}
\sum_{n \in \mathbb{N}} \eta_{n}<1 \tag{7.11}
\end{equation*}
$$

We will construct a sequence of mutually commuting finite type I factors $K_{i}$ and faithful states $\phi_{i}$ on $K_{i}$ such that

$$
\begin{equation*}
x_{j} \in K^{(n)}+V_{n}, \quad j=1, \ldots, n \tag{7.12}
\end{equation*}
$$

where

$$
\begin{equation*}
K^{(n)}=\left(\bigotimes_{i=1}^{n} K_{i}\right) \otimes 1 \tag{7.13}
\end{equation*}
$$

and

$$
\begin{equation*}
\|\left.\psi\right|_{\left(K^{(n-1)}\right)^{c}-\left.\phi_{n} \otimes \psi\right|_{\left(K^{(n)}\right)} \|<\eta_{n},} \tag{7.14}
\end{equation*}
$$

where $\left(K^{(n)}\right)^{c}=\left(K^{(n)}\right)^{\prime} \cap M$. It will then follow that $M \sim \bigotimes_{n \in \mathbb{N}}\left(K_{n}, \phi_{n}\right)$.
By assumption there is a finite type I factor $K_{1}$ and a faithful state $\phi_{1}$ such that eqs. (7.12) and (7.14) are satisfied for $n=1$ (use remark 7.2). Now assume that $K_{1}, \ldots, K_{n}$ and $\phi_{1}, \ldots, \phi_{n}$ have been chosen so that eqs. (7.12) and (7.14) are satisfied. Let $e_{i j}^{(n)}, i, j=1, \ldots, k_{n}$ be a complete set of matrix units for $K^{(n)}$, and let $e=e_{11}^{(n)}$. Since $M$ is properly infinite, $e \sim 1$ and $M_{e}=e M e \sim M$. In particular $M_{e}$ has the product property. Let $W$ be a strong neighbourhood of 0 such that the sum of any $k_{n}^{2}$ elements from $W$ must lie in $V_{n+1}$. There exists a finite type I subfactor $L_{n+1}$ of $M_{e}, y_{k, i, j}^{(n)} \in L_{n+1}, i, j=1, \ldots, k_{n}, k=1, \ldots, n+1$, and a state $\phi_{n+1}$ such that

$$
\begin{equation*}
e_{i 1}^{(n)} y_{k, i j}^{(n)} e_{1 j}^{(n)}-e_{i i}^{(n)} x_{k} e_{j j}^{(n)} \in W \tag{7.15}
\end{equation*}
$$

and eq. (7.14) is satisfied with $n$ replaced by $n+1$ (where we have used the canonical identification of $M_{e}$ with $\left.K^{(n) c}\right)$. Let

$$
\begin{equation*}
y_{k}^{(n)}=\sum_{i, j=1}^{k_{n}} e_{i 1}^{(n)} y_{k, i j}^{(n)} e_{1 j}^{(n)} \in K^{(n+1)} \tag{7.16}
\end{equation*}
$$

$k=1, \ldots, n+1$, where $K_{n+1}$ is the finite type I subfactor of $M$ obtained from the canonical embedding of $L_{n+1}$. Then

$$
y_{k}^{(n)}-x_{k} \in V_{n+1}
$$

so that eq. (7.12) is satisfied. Now let $\phi=\bigotimes_{k=1}^{\infty} \phi_{k}$. It follows from eqs. (7.11) and (7.14) that $\|\phi-\psi\|<1$. Hence the representation $\pi_{\phi}$ of the UHF $C^{*}$-algebra $\mathscr{A}=\bigcup_{j \in \mathbb{N}} K_{j}$ induced by $\phi$ is unitarily equivalent to $\pi_{\psi}$ (see for example [16, theorem 2.7]). Eq. (7.12) implies that $\pi_{\psi}(\mathscr{A})^{\prime \prime}=M$.

Remark 7.7. By further arguments, which we omit, one can show that in the general case the product property implies ITPFI.

## 8. ITPFI factors and AT flows

In this section we prove our major result, namely the equivalence of the ITPFI and AT properties (theorem 8.3). The proof that ITPFI implies AT (lemma 8.1) is rather straightforward. It follows from the known existence of a factor martingale (lemma
7.3), together with the relationship between certain positive operators in the centralizer and finite measures on the flow space (see eqs. (4.13)-(4.18) and (4.30)(4.33)). The converse is obtained by proving that AT of the flow of weights implies the product property (lemma 8.2).
Lemma 8.1. Let $M$ be an ITPFI factor of type $\mathrm{III}_{0}$. Then the flow of weights for $M$ is AT.
Proof. Since the proof for the discrete case $T(M) \neq\{0\}$ is much more transparent and, furthermore, motivates the proof for the general case, we present it first.
Case (i). $T(M) \neq\{0\}$. Let $T \in T(M), T \neq 0$. Using either lemma 11.2 of $[1]$ or results from [2] we can write $M=\bigotimes_{k=1}^{\infty}\left(M_{k}, \phi_{k}\right)$ where $\sigma_{T}^{\phi_{k}}=1$ for all $k$. We will use the discrete construction of the flow of weights here. Let $\lambda, \tilde{M}, \tilde{\phi}, \theta, B$ and $\nu_{B}$ be as in eqs. (4.1)-(4.12). We will prove that the base transformation $\theta$ is AT. By lemma 2.5 this is equivalent to the flow being AT.

Let $\varepsilon>0$, and let $\mu_{1}, \ldots, \mu_{n}$ be finite measures on the base space $B, \mu_{1}, \ldots, \mu_{n}<$ $\nu_{B}$. Using eqs. (4.15) and (4.18) we obtain projections $e_{j} \in \tilde{M}_{\tilde{\phi}}$ such that $\mu_{e_{j}}=\mu_{j}$, $j=1, \ldots, n$. It follows from our lemma 7.3 and lemma 2.3 of [14] that there exist $m<\infty$ and positive operators $f_{j} \in M_{\phi^{(m)}}^{(m)}$ where $M^{(m)}=\bigotimes_{k=0}^{m} M_{k}, \phi^{(m)}=\omega_{\lambda} \otimes$ $\left(\otimes_{k=1}^{m} \phi_{k}\right)$ and $M_{0}=\mathscr{L}\left(\ell^{2}(\mathbb{Z})\right)$, such that

$$
\begin{equation*}
\tilde{\phi}\left(\left|e_{j}-f_{j}\right|\right) \leq \varepsilon, \quad j=1, \ldots, n \tag{8.1}
\end{equation*}
$$

Since $e_{j}, f_{j} \in M_{\phi}$ we have the inequality

$$
\left|\tilde{\phi}\left(\left(e_{j}-f_{j}\right) c\right)\right| \leq\|c\| \tilde{\phi}\left(\left|e_{j}-f_{j}\right|\right), \quad c \in M
$$

(see for example, [7, lemma 11, p. 63]. It now follows from eqs. (8.1) and (4.15) that

$$
\begin{equation*}
\left\|\mu_{e_{j}}-\mu_{f_{i}}\right\| \leq \varepsilon, \quad j=1, \ldots, n . \tag{8.2}
\end{equation*}
$$

The proof will be completed by showing that

$$
\begin{equation*}
f_{j}=\sum_{k \in \mathbb{Z}} \sum_{l=1}^{N} c_{j k l} e_{j k l}, \quad j=1, \ldots, n, \tag{8.3}
\end{equation*}
$$

where $c_{j k l} \geq 0, e_{j k l}$ are minimal projections in $M_{\phi^{(m)}}^{(m)}, e_{j k l} \sim \theta^{k} e$ where $e$ is a fixed minimal projection in $M_{\phi^{(m)}}^{(m)}$, and $N<\infty$. Eqs. (8.3) and (4.15) and lemma 5.8 then give

$$
\begin{equation*}
\mu_{f_{j}}=\sum_{k \in \mathbb{Z}} c_{j k} \lambda^{-k} \theta^{k} \mu_{e}, \quad j=1, \ldots, n \tag{8.4}
\end{equation*}
$$

where $c_{j k}=\sum_{l} c_{j k t}$. The AT now follows directly from eqs. (8.2) and (8.4). Eq. (8.3) will follow directly from the structure of $M_{\phi^{(m)}}^{(m)}$ which we now determine.

Recall that if $\psi(x)=\operatorname{Trace}(\rho x), x \in P$ is a weight on a type I factor $P$ and $\rho \in P$, then $\sigma_{t}^{\psi}=\operatorname{Ad} \rho^{i t}$ and hence $P_{\psi}=\{\rho\}^{\prime}$. Thus the determination of the structure of $M_{\phi^{(m)}}^{(m)}$ becomes an elementary exercise in linear algebra. We have $\phi^{(m)}(x)=$ Trace $\left(\rho^{(m)} x\right), x \in M^{(m)}$ where $\rho^{(m)} \in M^{(m)}$ has eigenvalues $\prod_{k=0}^{m} \lambda_{k, j(k)}$ where $\left\{\lambda_{k}\right\}_{j=1, \ldots, n_{k}}$ is the eigenvalue list of $\left(M_{k}, \phi_{k}\right)$ and $n_{k}<\infty$ for $k \neq 0$. Since $\sigma_{T}^{\phi_{k}}=1$, all ratios $\lambda_{k r} / \lambda_{k s}$ are precisely some integral power of $\lambda$. For $k=0$, each $\lambda^{p}, p \in \mathbb{Z}$ occurs precisely once as an eigenvalue. Hence the eigenvalues of $\rho^{(m)}$ are of the
form $\alpha \lambda^{s}$ where $\alpha$ is fixed, and each $s \in \mathbb{Z}$ occurs precisely

$$
\begin{equation*}
N=\prod_{k=1}^{m} n_{k} \tag{8.5}
\end{equation*}
$$

times. It follows that

$$
\begin{equation*}
M_{\phi}^{(m)}=\left\{\rho^{(m)}\right\}^{\prime}=\oplus_{k \in \mathbb{Z}} M(k) \tag{8.6}
\end{equation*}
$$

where $M(k)=F_{N}$, and eqs. (4.1), (4.2) and (4.8) and $\theta=\operatorname{Ad} \tilde{S}$ imply that $\theta(M(k))=$ $M(k+1)$. Eq. (8.3) now results directly from diagonalizing $f_{j}$.
Case (ii) is the general case. We use here the continuous construction of the flow of weights. Let $\tilde{M}, \tilde{\omega}, \theta_{s}, X$ and $\nu$ be as in eqs. (4.19)-(4.29). Let $\varepsilon>0$, and let $\mu_{1}, \ldots, \mu_{n}$ be finite measures on the flow space $X, \mu_{1}, \ldots, \mu_{n}<\nu$. Precisely as in case (i) there exist $e_{j} \in \tilde{M}_{\bar{\omega}}$ such that $\mu_{e_{j}}=\mu_{j}$ (see eqs. (4.32), (4.33)), and positive operators $f_{j} \in M_{\phi}^{(m)}\left(\frac{m}{m}\right)$ for some $m<\infty$, satisfying eq. (8.2). The same argument used above shows that this time we have

$$
\begin{equation*}
M_{\phi^{(m)}}^{(m)}=\int^{\oplus} d t M(t) \tag{8.7}
\end{equation*}
$$

where $M(t)=F_{N}$ and $N$ is again given by eq. (8.5). Let $f \in M_{\phi}^{(m)}$. Then

$$
\begin{gather*}
f=\int^{\oplus} d t f(t)  \tag{8.8}\\
\theta_{s} f=\int^{\oplus} d t f(t-s) \tag{8.9}
\end{gather*}
$$

and

$$
\begin{equation*}
\tilde{\omega}(f)=\int d t e^{t} \tau(f(t)) \tag{8.10}
\end{equation*}
$$

where $\tau$ is the trace on $F_{N}$. Diagonalizing the positive operators $f_{j}$, we obtain

$$
\begin{equation*}
f_{j}=\sum_{l=1}^{N} f_{j l} \tag{8.11}
\end{equation*}
$$

where

$$
\begin{equation*}
f_{j t}=\int^{\oplus} d t f_{j l}(t) \tag{8.12}
\end{equation*}
$$

and

$$
\begin{equation*}
f_{j l}(t)=c_{j l}(t) e_{j l}(t) \tag{8.13}
\end{equation*}
$$

where $c_{j l}(t) \geq 0$ and $e_{j l}(t)$ are minimal projections in $M(t)$. Since $\tilde{\omega}\left(f_{j}\right)=\mu_{f_{j}}(x)<\infty$, eq. (8.10) implies that $e^{\prime} c_{j l}(t) \in L^{1}(\mathbb{R})$. Since the transitive action of $\mathbb{R}$ on $\mathbb{R}$ is AT, there exist $\lambda_{j l} \in L^{1}(\mathbb{R})$ and $c(t), e^{t} c(t) \in L^{1}(\mathbb{R})$ such that

$$
\begin{equation*}
\left\|e^{t} c_{j l}(t)-\int d s \lambda_{j l}(s) e^{s+t} c(s+t)\right\|_{1} \leq \varepsilon N^{-1} \tag{8.14}
\end{equation*}
$$

Let

$$
\begin{equation*}
F_{j l}=\int^{\oplus} d t c_{j l}(t) e(t) \tag{8.15}
\end{equation*}
$$

where $e(t)$ is a measurable family of minimal projections in $M(t)$. Since $e(t) \sim e_{j l}(t)$ in $M(t)$ we have $F_{j l} \sim f_{j l}$ in $M_{\phi^{\prime m}}^{(m)}$ and hence

$$
\begin{equation*}
\mu_{F_{j}}=\mu_{f_{j i}} \tag{8.16}
\end{equation*}
$$

Let

$$
\begin{equation*}
C=\int^{\oplus} d t c(t) e(t) \tag{8.17}
\end{equation*}
$$

and

$$
\begin{equation*}
G_{j l}=\int^{\oplus} d t g_{j l}(t) e(t) \tag{8.18}
\end{equation*}
$$

where

$$
\begin{equation*}
g_{j l}(t)=\int d s \lambda_{j l}(s) e^{s} c(s+t) \tag{8.19}
\end{equation*}
$$

Then

$$
\begin{equation*}
G_{j l}=\int d s \lambda_{j l}(s) e^{s} \theta_{s} C \tag{8.20}
\end{equation*}
$$

Lemma 4.1 now gives

$$
\begin{equation*}
\mu_{G_{j}}=\int d s \lambda_{j l}(s) e^{s} \theta_{s} \mu_{C} \tag{8.21}
\end{equation*}
$$

Using eqs. (8.10), (8.14), (8.15), (8.18) and (8.19) we obtain

$$
\begin{equation*}
\left\|\mu_{F_{j}}-\mu_{G_{j l}}\right\|=\int d t e^{t}\left|c_{j}(t)-g_{j l}(t)\right| \leq \varepsilon N^{-1} \tag{8.22}
\end{equation*}
$$

The AT of $\theta_{s}$ now follows from eqs. (8.2), (8.11), (8.16), (8.21) and (8.22).
Lemma 8.2. Let $M$ be a Krieger factor of type $\mathrm{III}_{0}$. If the flow of weights is AT , then $M$ has the product property.
Proof. Let $\varepsilon>0, x_{1}, \ldots, x_{p} \in M, V$ a strong neighbourhood of 0 in $M$, and $\phi$ a faithful state on $M$. We give first an outline of the argument. Connes' martingale condition [3] gives the existence of a conditional expectation $E$ onto a finitedimensional subalgebra $N$ such that $\phi \circ E=\phi$, and $x_{1}, \ldots, x_{p} \in N+V$. If $N$ were a finite type I subfactor, the condition $\phi \circ E=\phi$ would imply that $\phi$ was a product state relative to $M=N \otimes N^{c}$ and the product property would be trivially satisfied. The strategy is to embed $N$ in a finite type I subfactor $P$ so that $\phi$ is approximately a product state. The basic idea is to use the AT condition to construct a new state $\psi$ such that $\|\phi-\psi\| \leq \varepsilon$, and a set of matrix units for $P$ which are eigenvectors of $\sigma_{t}^{\psi}$ (which implies that $\psi$ is a product state relative to $M=P \otimes P^{c}$ ). More precisely, one selects a minimal projection $e^{(k)}$ from each full matrix algebra in $N$. A measure $\mu_{k}$ is associated with each $e^{(k)}$. Using the AT condition the $\mu_{k}$ can be approximated by measures $\mu_{k}^{\prime}=\sum \mu_{k l}$. The $\mu_{k l}$ determine both the desired subprojections of the $e^{(k)}$ (which are minimal projections in $P$ ), and the state $\psi$.

We give first the argument for the discrete case $T(M) \neq\{0\}$. The argument for the general case then proceeds in exactly the same way.
Case (i). There exists $T \in T(M), T \neq 0$. Let $\phi$ be a faithful state on $M$ such that $\sigma_{T}^{\phi}=1$.
Step (i). By [3] there is a conditional expectation $E: M \rightarrow N$ where $N$ is a finitedimensional subalgebra of $M$ such that

$$
\begin{equation*}
\phi(E(x))=\phi(x) \quad \text { for all } x \in M \tag{8.23}
\end{equation*}
$$

and

$$
\begin{equation*}
E\left(x_{j}\right) \in V+x_{j}, \quad j=1, \ldots, p \tag{8.24}
\end{equation*}
$$

We have

$$
\begin{equation*}
N=\oplus_{k=1}^{n} N_{k}, \tag{8.25}
\end{equation*}
$$

where the $N_{k}$ are type $I_{m_{k}}$ factors. Let $\phi_{k}$ denote the restriction of $\phi$ to $N_{k}$. Then

$$
\begin{equation*}
\phi_{k}(x)=\text { Trace } \rho_{k} x, \quad x \in N_{k}, \tag{8.26}
\end{equation*}
$$

where $\rho_{k} \in N_{k}$. Choose matrix units $e_{i j}^{k}$ for $N_{k}$ such that $\rho_{k}$ is diagonal, i.e.

$$
\begin{equation*}
\rho_{k} e_{i j}^{k}=\lambda_{k i} e_{i j}^{k} \delta_{i j}, \quad i, j=1, \ldots, m_{k} \tag{8.27}
\end{equation*}
$$

and

$$
\begin{equation*}
\lambda_{k 1} \geq \lambda_{k 2} \geq \cdots \geq \lambda_{k m_{k}}>0 \tag{8.28}
\end{equation*}
$$

Then, since $\sigma_{t}^{\phi \circ E}(E(x))=\sigma_{t}^{\phi}(E(x))$, we have

$$
\begin{equation*}
\sigma_{t}^{\phi}\left(e_{i j}^{k}\right)=\sigma_{t}^{\phi_{k}}\left(e_{i j}^{k}\right)=\left(\lambda_{k i} / \lambda_{k j}\right)^{i t} e_{i j}^{k} \tag{8.29}
\end{equation*}
$$

Step (ii). In order to use the AT condition, we construct the flow of weights as in $\S 4$, eqs. (4.1)-(4.18). Let $P_{0}$ denote the projection onto the basis vector $e_{0}$ of $\ell^{2}(\mathbb{Z})$. Let

$$
\begin{equation*}
f^{k}=P_{0} \otimes e^{k}, \quad k=1, \ldots, n \tag{8.30}
\end{equation*}
$$

where $e^{k}=e_{11}^{k}$. Since $P_{0} \in \mathscr{L}\left(\ell^{2}(\mathbb{Z})\right)_{\omega_{\lambda}}$ and $e^{k} \in M_{\phi}$ (see eq. (8.29)) we have $f^{k} \in \tilde{M}_{\bar{\phi}}$. Since $\tilde{\phi}\left(f^{k}\right)=\phi\left(e^{k}\right)<\infty$, eq. (4.15) defines measures

$$
\begin{equation*}
\mu_{k}=\mu_{f^{k}}, \quad k=1, \ldots, n, \tag{8.31}
\end{equation*}
$$

on the base space $B$ such that $\mu_{k}<\nu_{B}$. Using the AT condition and a variant of remark 2.2 (iii) we obtain a measure $\mu<\nu_{B}$ and integers $n_{k}(l), k=1, \ldots, n$, $l=-L, \ldots, L$ such that

$$
\begin{equation*}
\left\|\mu_{k}-\mu_{k}^{\prime}\right\| \leq \frac{1}{2} n^{-1} m_{k}^{-1} \varepsilon, \quad k=1, \ldots, n \tag{8.32}
\end{equation*}
$$

where the measures

$$
\begin{equation*}
\mu_{k}^{\prime}=\sum_{l=-L}^{+L} n_{k}(l) \lambda^{-l} \theta^{I} \mu \tag{8.33}
\end{equation*}
$$

are non-zero.
Step (iii). We now modify the state $\phi$. We first show that $\phi$ is determined by its restrictions to $M_{e^{k}}, k=1, \ldots, n$ where, in effect, it is the conditional expectation. For $x \in M$ we have

$$
\begin{equation*}
E(x)=\sum x_{i j}^{k} e_{i j}^{k} \tag{8.34}
\end{equation*}
$$

where the numbers $x_{i j}^{k}$ are determined by the equation

$$
\begin{equation*}
E\left(e_{i i}^{k} x e_{j j}^{k}\right)=x_{i j}^{k} e_{i j}^{k} \tag{8.35}
\end{equation*}
$$

Eqs. (8.23), (8.26), (8.27) and (8.34) give

$$
\begin{equation*}
\phi(x)=\sum \lambda_{k i} x_{i i}^{k} . \tag{8.36}
\end{equation*}
$$

Since

$$
\begin{equation*}
\phi\left(e_{1 i}^{k} x e_{i 1}^{k}\right)=\phi\left(e_{1 i}^{k} E(x) e_{i 1}^{k}\right)=\lambda_{k} x_{i i}^{k}, \tag{8.37}
\end{equation*}
$$

we get

$$
\begin{equation*}
\phi(x)=\sum_{k=1}^{n} \sum_{i=1}^{m_{k}}\left(\lambda_{k i} / \lambda_{k 1}\right) \phi\left(e_{1 i}^{k} x e_{i 1}^{k}\right) . \tag{8.38}
\end{equation*}
$$

The desired alteration of $\phi$ will now be obtained by changing it on $M_{e^{k}}$ and using eq. (8.38) to define the new state $\psi$ on $M$. From eq. (8.32) and lemma 5.9 we obtain states $\tilde{\psi}_{k}$ such that

$$
\begin{gather*}
s\left(\tilde{\psi_{k}}\right)=f^{k},  \tag{8.39}\\
\mu_{\tilde{\psi}_{k}}=\mu_{k}^{\prime}, \tag{8.40}
\end{gather*}
$$

and

$$
\begin{equation*}
\left\|\tilde{\psi}_{k}-\tilde{\phi}_{k}\right\| \leq n^{-1} m_{k}^{-1} \varepsilon, \tag{8.41}
\end{equation*}
$$

where $\tilde{\phi}_{k}$ is defined on $\tilde{M}_{f^{k}}$ by

$$
\begin{equation*}
\tilde{\phi}_{k}\left(P_{0} \otimes x\right)=\phi_{k}(x) . \tag{8.42}
\end{equation*}
$$

We define $\psi_{k}$ on $M_{e^{k}}$ by

$$
\begin{equation*}
\psi_{k}(x)=\tilde{\psi}_{k}\left(P_{0} \otimes x\right) \tag{8.43}
\end{equation*}
$$

We define $\psi$ on $M$ by

$$
\begin{equation*}
\psi(x)=\sum_{k=1}^{n} \sum_{i=1}^{m_{k}}\left(\lambda_{k j} / \lambda_{k 1}\right) \psi_{k}\left(e_{1 j}^{k} x e_{j 1}^{k}\right) . \tag{8.44}
\end{equation*}
$$

Eqs. (8.28) and (8.41)-(8.44) give

$$
\begin{equation*}
\|\psi-\phi\| \leq \varepsilon . \tag{8.45}
\end{equation*}
$$

Step (iv). We now embed $N$ in a type I factor $P \subset M$ so that $\psi$ is a product state relative to $M=P \otimes P^{c}$. Using eqs. (8.33), (8.40) and lemma 5.7 we obtain non-zero projections $f_{j}^{k}, j=1, \ldots, q_{k}, k=1, \ldots, n$ such that

$$
\begin{align*}
& f^{k}=\sum_{j=1}^{q_{k}} f_{j}^{k},  \tag{8.46}\\
& \tilde{\psi}_{k}=\sum_{j=1}^{q_{k}} \tilde{\psi}_{k j}, \tag{8.47}
\end{align*}
$$

where

$$
\begin{equation*}
\tilde{\psi}_{k j}=\left(\tilde{\psi}_{k}\right)_{f_{i}^{k}}, \tag{8.48}
\end{equation*}
$$

and integers $s_{k j}^{\prime}$ such that

$$
\begin{equation*}
\boldsymbol{\mu}_{\tilde{\psi}_{k j}}=\lambda^{-s_{k j}^{\prime}} \theta^{s_{k j}} \boldsymbol{\mu}=\lambda^{-s_{k, j}} \theta^{s_{k j}} \mu_{\tilde{\psi}_{11}}, \tag{8.49}
\end{equation*}
$$

where $s_{k j}=s_{k j}^{\prime}-s_{11}$. It now follows from lemmas 5.5 and 5.8 that

$$
\begin{equation*}
\tilde{\psi}_{k j} \sim \lambda^{s_{k j}} \tilde{\psi}_{11} . \tag{8.50}
\end{equation*}
$$

Since $P_{0}$ is one-dimensional, we can define projections $e_{j}^{k} \in M$ by

$$
\begin{equation*}
f_{j}^{k}=P_{0} \otimes e_{j}^{k} \tag{8.51}
\end{equation*}
$$

and it then follows from eqs. (8.44), (8.47), (8.48), (8.50), (8.51) and lemma 1.2.3(b) of [2] that there exist partial isometries $u_{j 1}^{k!} \in M$ such that

$$
\begin{align*}
& \left(u_{j 1}^{k 1}\right)^{*} u_{j 1}^{k 1}=e_{1}^{1},  \tag{8.52}\\
& u_{j 1}^{k 1}\left(u_{j 1}^{k k}\right)^{*}=e_{j}^{k}, \tag{8.53}
\end{align*}
$$

and

$$
\begin{equation*}
\sigma_{t}^{\psi}\left(u_{j 1}^{k 1}\right)=\lambda^{i t s_{k j}} u_{j 1}^{k 1}, \tag{8.54}
\end{equation*}
$$

$j=1, \ldots, q_{k}, k=1, \ldots, n$. We extend the range of the index $j$ to $1, \ldots, m_{k} q_{k}$ by defining

$$
\begin{equation*}
u_{j 1}^{k 1}=e_{11}^{k} u_{11}^{k!} \tag{8.55}
\end{equation*}
$$

where $j=(l-1) q_{k}+t, \mathrm{I} \leq t<q_{k}, l=1, \ldots, m_{k}$. Finally we obtain a complete set of matrix units by defining

$$
\begin{equation*}
u_{i j}^{k l}=u_{i 1}^{k 1}\left(u_{j 1}^{l 1}\right)^{*} \tag{8.56}
\end{equation*}
$$

where $i=1, \ldots, q_{k}, j=1, \ldots, q_{l}, k, l=1, \ldots, n$. We can now define $P$ to be the type I factor generated by the $u_{i j}^{k l}$. Eqs. (8.44) and (8.54)-(8.56) imply that the $u_{i j}^{k l}$ are eigenvectors of $\sigma_{t}^{\psi}$. Hence $\sigma_{t}^{\psi}(P)=P$ which implies that $\psi$ is a product state relative to $M=P \otimes P^{c}$.
Case (ii) is the general case. The argument is virtually identical. We use here the lacunary construction of the flow of weights. For this purpose we take $\phi$ to be a faithful lacunary integrable weight of infinite multiplicity. Basically the only change is that lemmas $5.5,5.7,5.8$ and 5.9 are replaced by lemmas $6.1,6.2$ and 6.4 .

Step (i) is identical. In step (ii) we now construct the flow of weights as in $\S 4$, eqs. (4.39)-(4.48). In particular eq. (4.47) defines measures

$$
\begin{equation*}
\mu_{k}=\mu_{\phi_{e_{k}}}, \quad k=1, \ldots, n \tag{8.57}
\end{equation*}
$$

on the flow space $X, \mu_{k}<\nu$. Using the AT condition and again a variation on remark 2.2 (iii) we obtain a measure $\mu<\nu, t_{l} \in \mathbb{R}_{+}^{*}, n_{k}(l) \in \mathbb{Z}, l=1, \ldots, L$ such that

$$
\begin{equation*}
\left\|\mu_{k}-\mu_{k}^{\prime}\right\| \leq \frac{1}{5} n^{-1} m_{k}^{-1} \varepsilon, \quad k=1, \ldots, n, \tag{8.58}
\end{equation*}
$$

where the measures

$$
\begin{equation*}
\mu_{k}^{\prime}=\sum_{t=1}^{L} n_{k}(l) t^{-1} \mathscr{F}_{t_{t}}^{M} \mu \tag{8.59}
\end{equation*}
$$

are non-zero.
Step (iii) is almost identical. We work directly on $M$ and use lemma 6.4 to obtain states $\psi_{k}$ on $M_{e^{k}}$ satisfying

$$
\begin{align*}
s\left(\psi_{k}\right) & =e^{k},  \tag{8.60}\\
\mu_{\psi_{k}} & =\mu_{k}^{\prime}, \tag{8.61}
\end{align*}
$$

and

$$
\begin{equation*}
\left\|\psi_{k}-\phi_{k}\right\| \leq n^{-1} m_{k}^{-1} \varepsilon . \tag{8.62}
\end{equation*}
$$

$\psi$ is again defined by eq. (8.44) and satisfies eq. (8.45). In step (iv) we begin by noting that the states $\psi_{k}$ defined by eq. (8.60) are integrable since the measures $\mu_{k}^{\prime}$ are smooth (see eq. (4.48) et seq.). Using eqs. (8.59), (8.61) and lemma 6.2 we obtain non-zero projections $e_{j}^{k}, j=1, \ldots, q_{k}, k=1, \ldots, n$ such that

$$
\begin{equation*}
e^{k}=\sum_{j=1}^{q_{k}} e_{j}^{k} \tag{8.63}
\end{equation*}
$$

and

$$
\begin{equation*}
\psi_{k}=\sum_{j=1}^{q_{k}} \psi_{k j}, \tag{8.64}
\end{equation*}
$$

where $\psi_{k_{j}}=\left(\psi_{k}\right)_{e_{j}^{k}}$, and $t_{k j}^{\prime} \in \mathbb{R}$ such that
where $t_{k j}=t_{k j}^{\prime}-t_{11}$. It now follows from lemma 6.1 that

$$
\begin{equation*}
\psi_{k j} \sim e^{-t_{k j}} \psi_{11} \tag{8.66}
\end{equation*}
$$

As before, it follows from eqs. (8.44), (8.66) and lemma 1.2.3(b) of [2] that there exist partial isometries $u_{i j}^{k l}$ satisfying eqs. (8.52)-(8.56), and the type I factor $P$ generated by the $u_{i j}^{k l}$ has the desired properties.

Theorem 8.3. Let $M$ be a type $\mathrm{III}_{0}$ injective factor. Then the following are equivalent.
(i) $M$ is ITPFI.
(ii) The flow of weights for $M$ is approximately transitive.
(iii) $M$ has the product property.

Proof. By [4] $M$ is a Krieger factor. We have the following implications: (i) $\Rightarrow$ (iii), corollary 7.4 ; (iii) $\Rightarrow$ (i), lemma 7.6 ; (i) $\Rightarrow$ (ii), lemma 8.1; and (ii) $\Rightarrow$ (iii), lemma 8.2.

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