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| Title | The one cocycle property for shifts |
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| Author(s) | KISHIMOTO, A. |
| Citation | Ergodic Theory and Dynamical Systems, 25, 823-859 <br> https://doi.org/10.1017/S0143385704000860 |
| Issue Date | 2005 |
| Doc URL | Copyright © 2005 Cambridge University Press |
| Rights | article |
| Type | ETDS25.pdf |
| File Information |  |

Instructions for use

# The one-cocycle property for shifts 

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(Received 11 December 2003 and accepted in final form 21 September 2004)


#### Abstract

The two-sided shift on the infinite tensor product of copies of the $n \times n$ matrix algebra has the so-called Rohlin property, which entails the one-cocycle property, useful in analyzing cocycle-conjugacy classes. In the case $n=2$, the restriction of the shift to the gauge-invariant CAR algebra also has the one-cocycle property. We extend the latter result to an arbitrary $n \geq 2$. As a corollary it follows that the flow $\alpha$ on the Cuntz algebra $\mathcal{O}_{n}=C^{*}\left(s_{0}, s_{1}, \ldots, s_{n-1}\right)$ defined by $\alpha_{t}\left(s_{j}\right)=e^{i p_{j} t} s_{j}$ has the Rohlin property (for flows) if and only if $p_{0}, \ldots, p_{n-1}$ generate $\mathbb{R}$ as a closed sub-semigroup. Note that then such flows are all cocycle-conjugate to each other.


## 1. Introduction

For an integer $n$ greater than 1 we denote by $M_{n}$ the $n \times n$ matrix algebra over the complex numbers $\mathbb{C}$. For each integer $m \in \mathbb{Z}$ we assign a copy $M_{n}^{(m)}$ of $M_{n}$ and take the infinite tensor product $B_{n}=\bigotimes_{m \in \mathbb{Z}} M_{n}^{(m)}$. The shift automorphism $\sigma$ of $B_{n}$ is defined by sending an element of $M_{n}^{(m)}$ to the corresponding element in $M_{n}^{(m+1)}$.

In [4] it was shown that $\sigma$ has the Rohlin property (see below) in the case $n=2$. The proof is based on a known connection between such a $\sigma$ and a certain quasi-free automorphism of the $C^{*}$-algebra associated with the canonical anti-commutation relations, or the CAR algebra.

This prompts us to attempt to generalize it. A further exploitation of the CAR algebra formalism was done in [2, 3]. A full generalization for any $n \geq 2$ was done in [15, 16]. Also, an extension to the shift on $B=\bigotimes_{\mathbb{Z}}\left(M_{2} \oplus M_{3}\right)$ was done in [18]. In this generality the Rohlin property for the shift $\sigma$ reads as follows: for any $N \in \mathbb{N}$ and $\epsilon>0$ there is a family $\left\{e_{1 i} \mid i=0,1, \ldots, N-1\right\} \cup\left\{e_{2 i} \mid i=0,1, \ldots, N\right\}$ of projections in $B$ such that $\sum_{i} e_{1 i}+\sum_{2 i} e_{2 i}=1$ and

$$
\left\|\sigma\left(e_{1, i}\right)-e_{1, i+1}\right\|<\epsilon
$$

for $i=0, \ldots, N-2$ and

$$
\left\|\sigma\left(e_{2, i}\right)-e_{2, i+1}\right\|<\epsilon
$$

for $i=0,1, \ldots, N-1$. (Hence, it follows that $\left.\sigma\left(e_{1, N-1}+e_{2, N}\right) \approx e_{1,0}+e_{2,0}.\right)$

A useful consequence of this property is the one-cocycle property, i.e. for any unitary $u$ and $\epsilon>0$ there is a unitary $v$ such that $\left\|u-v \sigma(v)^{*}\right\|<\epsilon$, which is the property we actually need for applications. We refer the reader to $[\mathbf{7 , 1 0}]$ for the Rohlin property and more applications.

There is yet another attempt on the restriction of the shift to a certain $C^{*}$-subalgebra of $\otimes_{\mathbb{Z}} M_{2}$, which corresponds to the gauge-invariant CAR algebra, based on the CAR algebra formalism [19]. In this case the shift (on this subalgebra) cannot have the Rohlin property, but can have an approximate version of Rohlin property. The result is that the shift has the one-cocycle property (appropriately formulated).

In this note we extend the above result to the case of general $n>2$. We now formulate the problem more precisely below.

Let $\mathcal{U}_{n}$ be the unitary group of $M_{n}$ and define an action $\beta$ of $\mathcal{U}_{n}$ on $B_{n}=\bigotimes_{m \in \mathbb{Z}} M_{n}^{(m)}$ by

$$
\beta_{u}=\bigotimes_{m \in \mathbb{Z}} \operatorname{Ad} u^{(m)}
$$

where $u^{(m)}$ is the copy of $u$ in $M_{n}^{(m)}$. Note that $\beta_{u}$ commutes with $\sigma$.
Denoting by $\mathbb{T}$ the group of complex numbers of modulus one, we regard $\mathbb{T}^{n-1}$ as a subgroup of $\mathcal{U}_{n}$ by

$$
\left(z_{1}, \ldots, z_{n-1}\right) \mapsto\left(\begin{array}{cccc}
1 & & & \\
& z_{1} & & \\
& & \ddots & \\
& & & z_{n-1}
\end{array}\right)
$$

and define an action $\gamma$ of $\mathbb{T}^{n-1}$ by $\beta \mid \mathbb{T}^{n-1}$. We call $\gamma$ the gauge action (of $\mathbb{T}^{n-1}$ ) on $B_{n}$.
Let $\left\{e_{i j}\right\}$ denote the matrix units for $M_{n}$ and let

$$
v_{m}=\sum_{i j} e_{i j}^{(m)} e_{j i}^{(m+1)}
$$

Then it follows that $v_{m}$ is a self-adjoint unitary in the fixed-point algebra $B_{n}^{\beta}$ and satisfies Ad $v_{m}\left(x^{(m)} y^{(m+1)}\right)=y^{(m)} x^{(m+1)}$ for $x, y \in M_{n}$. If we set $V_{m}=v_{-m} v_{-m+1} \cdots v_{0} \cdots v_{m}$, then $V_{m} \in B_{n}^{\beta}$ satisfies that $\sigma(x)=\lim \operatorname{Ad} V_{m}(x)$ for $x \in B_{n}$.

We set $A_{n}=B_{n}^{\gamma}$, the fixed-point algebra of $B_{n}$ under $\gamma=\beta \mid \mathbb{T}^{n-1}$. Note that $\sigma$ restricts to $A_{n}$, which is denoted by $\sigma \mid A_{n}$ or simply $\sigma$. Our purpose is to prove that $\sigma \mid A_{n}$ has the one-cocycle property for all $n$.

To present a precise statement we should note that the approximately finite-dimensional $C^{*}$-algebra $A_{n}$ is prime and not simple. It has $n$ maximal ideals of codimension one.

For $i=0,1, \ldots, n-1$, let $I_{i}$ be the (closed, two-sided) ideal of $A_{n}$ generated by $e_{j j}^{(m)}$, $j \neq i, m \in \mathbb{Z}$. We then note that $\left(\bigotimes_{-M}^{M} M_{n}\right) \cap I_{i}$ is orthogonal to $\Pi_{-M}^{M} e_{i i}^{(m)}$. It follows that the quotient $A_{n} / I_{i}$ is isomorphic to $\mathbb{C}$; let $\varphi_{i}$ be the corresponding character on $A_{n}$.

THEOREM 1.1. Let $n=2,3, \ldots$ and let $B_{n}, A_{n}, \sigma, \varphi_{i}$ be as above. Then the shift $\sigma$ on $A_{n}$ has the following properties.
(1) For any unitary $u \in A_{n}$ such that $\varphi_{i}(u)=1$ for $i=0, \ldots, n-1$, there is a sequence $\left(v_{k}\right)$ of unitaries in $A_{n}$ such that $\left\|u-v_{k} \sigma\left(v_{k}\right)^{*}\right\| \rightarrow 0$.
(2) If $\left(u_{k}\right)$ is a sequence in the unitary group $\mathcal{U}\left(A_{n}\right)$ such that $\varphi_{i}\left(u_{k}\right)=1$ for $i=$ $0, \ldots, n-1$ and $\left\|\left[u_{k}, x\right]\right\| \rightarrow 0$ for all $x \in B_{n}$, then there is a sequence $\left(v_{k}\right)$ in $\mathcal{U}\left(A_{n}\right)$ such that $\left\|u_{k}-v_{k} \sigma\left(v_{k}\right)\right\| \rightarrow 0$ and $\left\|\left[v_{k}, x\right]\right\| \rightarrow 0$ for all $x \in B_{n}$.
(3) If $\left(u_{k}\right)$ is a sequence in $\mathcal{U}\left(A_{n}\right)$ such that $\varphi_{i}\left(u_{k}\right)=1$ for $i=0, \ldots, n-1$ and $\left\|\left[u_{n}, x\right]\right\| \rightarrow 0$ for all $x \in A_{n}$, then there is a sequence $\left(v_{k}\right)$ in $\mathcal{U}\left(A_{n}\right)$ such that $\left\|u_{k}-v_{k} \sigma\left(v_{k}\right)\right\| \rightarrow 0$ and $\left\|\left[v_{k}, x\right]\right\| \rightarrow 0$ for all $x \in A_{n}$.

The proof of this theorem occupies the following sections.
Note that $A_{n}$ is an approximately finite-dimensional $C^{*}$-algebra, i.e. the closure of the union of an increasing sequence of finite-dimensional $C^{*}$-algebras [1]. Since we have to compare projections in such a $C^{*}$-algebra, we need a nice description of the dimension group $K_{0}\left(A_{n}\right)$ of $A_{n}$, i.e. the ordered group generated by equivalence classes of projections [6]. If $n=2, A_{2}$ is already presented in [1], in terms of Bratteli diagram, and $K_{0}\left(A_{2}\right)$ is identified, by Renault [24], with the integer-coefficient polynomials $\mathbb{Z}[x]$, with the positivity defined by strict positivity on the open interval $(0,1)$, which was good enough in the discussions in [19]. In general, Handelman [9] shows that $K_{0}\left(A_{n}\right)$ is identified with $\mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ with the order defined as follows: $p$ is positive if $p$ is expressed as

$$
p=\sum_{|v|=K} c_{v}\left(1-x_{1}-\cdots-x_{n-1}\right)^{v_{0}} x_{1}^{v_{1}} \cdots x_{n-1}^{v_{n-1}}
$$

with non-negative coefficients $c_{v} \geq 0$, where the sum is taken over all $v \in \mathbb{Z}_{+}^{n}$ with $|v|=$ $\sum_{i=0}^{n-1} v_{i}=K$. However, there seem to be no clear criteria for positivity as the Renault's result for $n=2$, because the strict positivity on the interior of the ( $n-1$ )-simplex $\Lambda_{n}=\left\{\left(x_{1}, \ldots, x_{n-1}\right) \mid x_{i}>0, \sum_{i} x_{i} \leq 1\right\}$ with some boundary conditions is not enough.

By introducing a notion of vanguard for $p \in \mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ we give a sufficient condition for positivity, which is useful enough in the following arguments. If $p$ is expressed as above and $i=0,1, \ldots, n-1$, the $i$-vanguard of $p$ on the level $K$ is a certain subset of $v$ with $|v|=K$ such that $c_{v} \neq 0$ and there is no $w$ between $v$ and $K e_{i}$ with $c_{w} \neq 0$. We see that the $i$-vanguard of $p$ is essentially independent of $K$ as well as the coefficients $c_{v}$; hence, if $p$ is positive, those coefficients $c_{v}$ must be positive, which does not follow from the positivity of $p$ as a function on $\Lambda_{n}$. The sufficient condition states in Proposition 2.5 that if $p$ is strictly positive on $\Lambda_{n}$ possibly except for the vertices and the above-mentioned condition is satisfied for all the vanguards, then $p$ is positive as an element of $\mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$.

Note that $A_{n}$ has quotients which are of type I. We consider irreducible representations of $A_{n}$ whose images contain the compact operators, which are necessarily $\sigma$-covariant since $\sigma$ is approximately inner. In Lemma 4.4 we show that the shift on a certain quotient of $A_{n}$ of type I has the properties as in the theorem. This gives the first step of the proof by induction on $n$.

Then we show how to locally embed $A_{n}$ into $A_{n+1}$, almost intertwining the shifts. This process will show, with the known result for $n=2$ shown by the CAR algebra formalism (4.3), that the shift on $A_{n}$ has the approximate Rohlin property for all $n$ (see Lemma 4.6); by this property (Definition 4.2) we essentially mean that there is a sequence $\left(e_{i}^{(k)}\right)$ in the orthogonal family of $N$ projections in $A_{n}$ for arbitrarily large $N$ such that $\left\|\sigma\left(e_{i}^{(k)}\right)-e_{i+1}^{(k)}\right\| \rightarrow 0$ as $k \rightarrow \infty$ and $N\left[e_{0}^{(k)}\right](x)$ converges to 1 uniformly
in $x$ on every compact subset of $\Lambda_{n}$ except for the vertices, as well as the lower estimate $\left[e_{0}^{(k)}\right](x) \geq C\left(x_{0} x_{1} \cdots x_{n-1}\right)^{S}$ for some $S \in \mathbb{N}$ and $C>0$ independent of $k$, where $x_{0}=1-x_{1}-\cdots-x_{n-1}$, and some estimates on the vanguards of $\left[e_{0}^{(k)}\right]$ independent of $k$. We see that the additional conditions on $\left[e_{0}^{(k)}\right]$ play an important role in comparison with other projections.

With some explicit estimates, in $K_{0}\left(A_{n}\right)$, of the projections involved, we then proceed just as in [19]. The crucial induction step is given in Lemma 4.12.

In the above theorem the first property is weaker than the third.
Remark 1.2. Let $\alpha$ be an approximately inner automorphism of $A_{n}$. If $\alpha$ has the third property of Theorem 1.1, then it also has the first property. This follows because then $\alpha$ is almost conjugate to the shift $\sigma$, i.e. there is a sequence $\left(\phi_{m}\right)$ of approximately inner automorphisms of $A_{n}$ such that $\left\|\alpha-\phi_{m} \sigma \phi_{m}^{-1}\right\| \rightarrow 0$, as follows from Theorem 4.1 of [8]. If $u \in \mathcal{U}\left(A_{n}\right)$ such that $\varphi_{i}(u)=1$, we have a sequence $\left(v_{m, k}\right)$ in $\mathcal{U}\left(A_{n}\right)$ for each $m \in \mathbb{N}$ such that $\left\|\phi_{m}(u)-v_{m, k} \sigma\left(v_{m, k}\right)^{*}\right\| \rightarrow 0$ as $k \rightarrow \infty$. Then we can choose a sequence $\left(v_{m}\right)$ from $\left(\phi^{-1}\left(v_{m, k}\right)\right)$ such that $\left\|u-v_{m} \alpha\left(v_{m}\right)^{*}\right\| \rightarrow 0$. (Since $\phi_{m}$ may not extend to an automorphism of $B_{n}$, this method does not show that $\alpha$ also has the second property. If $\alpha$ does not extend to an automorphism of $B_{n}, \alpha$ is unlikely to have the second property.)

The following is a corollary of Theorem 1.1(2).
COROLLARY 1.3. Let $n, B_{n}, A_{n}, \sigma$ be as above. Then the shift $\sigma$ on $B_{n}$ has a sequence $\left(U_{m}\right)$ in $\mathcal{U}\left(A_{n}\right)$ such that $\sigma(x)=\lim _{m \rightarrow \infty} \operatorname{Ad} U_{m}(x)$ for $x \in B_{n}$ and $\lim _{m \rightarrow \infty} \| \sigma\left(U_{m}\right)-$ $U_{m} \|=0$.

Proof. We have defined $V_{m}=v_{-m} v_{-m+1} \cdots v_{0} \cdots v_{m} \in A_{n}$, which satisfies

$$
\sigma(x)=\lim \operatorname{Ad} V_{m}(x), \quad x \in B_{n}
$$

but does not satisfy $\lim \left\|\sigma\left(V_{m}\right)-V_{m}\right\|=0$. (By calculation, it follows that the spectrum of $V_{m} \sigma\left(V_{m}^{*}\right)$ is $\left\{1, e^{2 \pi i / 3}, e^{-2 \pi i / 3}\right\}$ and hence that $\left\|\sigma\left(V_{m}\right)-V_{m}\right\|=\sqrt{3}$ for $m \geq 1$.) Since $V_{m} \sigma\left(V_{m}\right)^{*} \in A_{n}, \varphi_{i}\left(V_{m} \sigma\left(V_{m}^{*}\right)\right)=1$, and $\left(V_{m} \sigma\left(V_{m}^{*}\right)\right)$ is a central sequence in $\mathcal{U}\left(B_{n}\right)$, the previous theorem gives us a sequence $\left(W_{m}\right)$ in $\mathcal{U}\left(A_{n}\right)$ such that $\| V_{m} \sigma\left(V_{m}^{*}\right)-$ $W_{m} \sigma\left(W_{m}^{*}\right) \| \rightarrow 0$ and $\left\|\left[W_{m}, x\right]\right\| \rightarrow 0$ for $x \in B_{n}$. Let $U_{m}=W_{m}^{*} V_{m} \in \mathcal{U}\left(A_{n}\right)$. Then $\left(U_{m}\right)$ satisfies the required conditions.

In connection with the above corollary a question arises naturally of whether ( $U_{m}$ ) can be chosen from the smaller fixed-point algebra $B_{n}^{\beta}$.

When $\ell$ is a positive integer, we denote by $\sigma^{(\ell)}$ the cyclic shift on $\bigotimes_{m=0}^{\ell-1} M_{n}^{(m)}$ defined by $x^{(m)} \mapsto x^{(m+1)}$ with $x^{(\ell)}=x^{(0)}$. We also denote by $\gamma^{(\ell)}$ the action of $\mathbb{T}^{n-1}$ on $\bigotimes_{m=0}^{\ell-1} M_{n}^{(m)}$ by $\bigotimes_{m=0}^{\ell-1} \operatorname{Ad} z^{(m)}, z \in \mathbb{T}^{n-1}$.

COROLLARY 1.4. For the shift $\sigma$ on $B_{n}$ and $\epsilon>0$ there exists an increasing sequence $\left(\ell_{k}\right)$ in $\mathbb{N}$, an automorphism $\phi$ of $B_{n}$, and a unitary $w \in A_{n}$ such that $\phi \gamma_{z}=\gamma_{z} \phi$ for $z \in \mathbb{T}^{n-1},\|w-1\|<\epsilon$, and the pair of $\operatorname{Ad} w \phi \sigma \phi^{-1}$ and $\gamma$ is conjugate to the pair of $\bigotimes_{k=1}^{\infty} \sigma^{\ell_{k}}$ and $\bigotimes_{k=0}^{\infty} \gamma^{\left(\ell_{k}\right)}$ on $\bigotimes_{k=1}^{\infty}\left(\bigotimes_{m=0}^{\ell_{k}-1} M_{n}^{(m)}\right) \cong B_{n}$.

Proof. For $a, b \in \mathbb{N}$ with $a<b$ we define

$$
u^{(a, b)}=\sum_{i, j} e_{i j}^{(a)} e_{j i}^{(b)} .
$$

Then $\operatorname{Ad} u^{(a, b)} \sigma$ acts on $\bigotimes_{m=a}^{b-1} M_{n}^{(m)}$ as the cyclic shift.
By using Theorem 1.1 we can derive the following: for any finite subset $\mathcal{F}$ of $\mathbb{Z}$ and $\epsilon>0$ there is a finite subset $\mathcal{G}$ of $\mathbb{Z}$ such that for any unitary $u \in A_{n} \cap \bigotimes_{m \notin \mathcal{G}} M_{n}^{(m)}$ there is a unitary $v \in A_{n} \cap \bigotimes_{m \notin \mathcal{F}} M_{n}^{(m)}$ such that $\left\|u-v \sigma\left(v^{*}\right)\right\|<\epsilon$.

By using this fact we define inductively a decreasing sequence ( $a_{k}$ ), an increasing sequence $\left(b_{k}\right)$, and sequences $\left(v_{k}\right)$ and $\left(w_{k}\right)$ of unitaries in $A_{n}$ such that $a_{1}=0$, $b_{1}=1, w_{k} v_{k} \sigma\left(v_{k}\right)^{*}=u^{\left(a_{k}, b_{k}\right)},\left\|w_{k}-1\right\|<2^{-k} \epsilon$, and $v_{k}, w_{k} \in \bigotimes_{Z_{k}} M_{n}$ for mutually disjoint family $\left(Z_{k}\right)$ of finite subsets of $\mathbb{Z}$. Then $\phi=\lim _{k} \operatorname{Ad}\left(v_{1} v_{2} \cdots v_{k}\right)$ is well defined as an automorphism of $B_{n}$, as well as $w=\lim _{k} w_{1} w_{2} \cdots w_{k}$ as a unitary in $A_{n}$. Then we check that $\operatorname{Ad} w \phi \sigma \phi^{-1}$ leaves $\bigotimes_{m=a_{k}}^{b_{k}-1} M_{n}$ invariant and acts as a cyclic shift on $\bigotimes_{m=a_{k}}^{a_{k-1}-1} M_{n} \otimes \bigotimes_{m=b_{k-1}}^{b_{k}-1} M_{n}$ for each $k \geq 1$ with $a_{0}=b_{0}=0$. Then with $\ell_{k}=b_{k}-b_{k-1}-a_{k}+a_{k-1}$ we can conclude the proof.

The above corollary shows that the action of $\mathbb{T}^{n-1} \times \mathbb{Z}$ on $B_{n}$ defined by $(z, k) \mapsto \gamma_{z} \sigma^{k}$, which does not leave any finite-dimensional $C^{*}$-subalgebra of $B_{n}$ invariant, is cocycle conjugate to an action which is of product type.

Remark 1.5. Let $z \in \mathbb{T}^{n-1}$ be such that $z^{m} \neq 1$ for any non-zero $m \in \mathbb{Z}$ and define an action $\mathbb{Z}^{2}$ on $B_{n}$ by $\alpha_{(a, b)}=\gamma_{z}^{a} \sigma^{b}$. Then $\alpha$ is cocycle conjugate to the action $\alpha^{\prime}$ defined by

$$
\alpha_{(a, b)}^{\prime}=\phi \circ\left(\gamma_{z}\right)^{a}\left(\operatorname{Ad} \phi^{-1}(w) \sigma\right)^{b} \circ \phi^{-1}=\gamma_{z}^{a}\left(\operatorname{Ad} w \phi \sigma \phi^{-1}\right)^{b}
$$

in the notation of the above corollary, which is of product type. We should note that such an action with the Rohlin property (i.e. each $\alpha_{(a, b)}$ has the Rohlin property for $\left.(a, b) \neq 0\right)$ is unique up to cocycle-conjugacy [23]. Hence, since $\alpha$ has the Rohlin property, $\alpha$ is cocycle conjugate to, e.g., the action $\alpha^{\prime \prime}$ defined by

$$
\alpha_{(a, b)}^{\prime \prime}=\gamma_{z_{1}}^{a} \gamma_{z_{2}}^{b},
$$

where $z_{1}, z_{2} \in \mathbb{T}^{n-1}$ is any pair which generates a copy of $\mathbb{Z}^{2}$ in $\mathbb{T}^{n-1}$. It seems that we still do not know if there is an action of $\mathbb{Z}^{2}$ on $B_{n}$ with the Rohlin property which is not cocycle conjugate to an action of product type.

Another corollary extends what is stated in [19] for $n=2$.
COROLLARY 1.6. Let $n$ be an integer greater than 1 and $\mathcal{O}_{n}$ be the Cuntz algebra generated by $n$ isometries $s_{0}, s_{1}, \ldots, s_{n-1}$ satisfying the relation $\sum_{j=0}^{n-1} s_{j} s_{j}^{*}=1$ (see [5]). For $p=\left(p_{0}, p_{1}, \ldots, p_{n-1}\right) \in \mathbb{R}^{n}$ define a flow $\alpha$ on $\mathcal{O}_{n}$ by $\alpha_{t}\left(s_{j}\right)=e^{i p_{j} t} s_{j}$ for $j=0, \ldots, n-1$. Then the following conditions are equivalent:
(1) $\left\{p_{0}, \ldots, p_{n-1}\right\}$ generates $\mathbb{R}$ as a closed sub-semigroup;
(2) the crossed product $\mathcal{O}_{n} \times_{\alpha} \mathbb{R}$ is purely infinite and simple;
(3) $\alpha$ has the Rohlin property, i.e. for any $\lambda \in \mathbb{R}$ there is a central sequence $\left(u_{m}\right)$ in $\mathcal{U}\left(\mathcal{O}_{n}\right)$ such that $\left\|\alpha_{t}\left(u_{m}\right)-e^{i \lambda t} u_{m}\right\| \rightarrow 0$ uniformly in $t$ on every compact subset of $\mathbb{R}$.

Proof. First we note that $\mathcal{O}_{n}$ is a purely infinite, simple, nuclear $C^{*}$-algebra [5] and such a class is now well studied $[\mathbf{1 2 , 1 3 ]}$.

The equivalence of (1) and (2) follows from [14, 22]. That (3) implies (2) follows from [17]. What is left to prove is that (1) implies (3). If (1) holds, it is shown by combinatorial arguments in [18] that for any $\lambda \in \mathbb{R}$ there is a sequence $\left(u_{m}\right)$ of unitaries in the ${ }^{*}$-subalgebra generated by $s_{0}, \ldots, s_{n-1}$ such that $\left\|\alpha_{t}\left(u_{m}\right)-e^{i \lambda t} u_{m}\right\| \rightarrow 0$ uniformly in $t$ on every compact subset of $\mathbb{R}$. As in the proof of Proposition 3.2 of [19], with the property of $\sigma \mid A_{n}$ as given in Theorem 1.1(1), there is a sequence $\left(\phi_{k}\right)$ of unital endomorphisms of $\mathcal{O}_{n}$ such that $\left[\phi_{k}(x), y\right] \rightarrow 0$ for any $x, y \in \mathcal{O}_{n}$ and $\alpha_{t} \phi_{k}=\phi_{k} \alpha_{t}$. Then we can choose a central sequence from $\left\{\phi_{k}\left(u_{m}\right) \mid k, m \in \mathbb{N}\right\}$ which satisfies the required condition for $\lambda \in \mathbb{R}$.

We would like to add that the flows $\alpha$ on $\mathcal{O}_{n}$ satisfying the conditions in the above corollary are cocycle conjugate to each other [19], i.e. for any flows $\alpha$ and $\alpha^{\prime}$ on $\mathcal{O}_{n}$ of the above form there is an automorphism $\phi$ of $\mathcal{O}_{n}$ and an $\alpha$-cocycle $u$ such that $\operatorname{Ad} u_{t} \alpha_{t}=\phi \alpha^{\prime} \phi^{-1}$. See $[\mathbf{2 0}, \mathbf{2 1}]$ for more on Rohlin flows.

## 2. The dimension group of $A_{n}$

For $n=2,3, \ldots$ we denote by $\mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ the abelian group of integer coefficient polynomials in $x_{1}, \ldots, x_{n-1}$. For $v=\left(v_{0}, v_{1}, \ldots, v_{n-1}\right) \in \mathbb{Z}_{+}^{n}$ we let

$$
x^{v}=x_{0}^{v_{0}} x_{1}^{v_{1}} \cdots x_{n-1}^{v_{n-1}},
$$

where $x_{0}=1-x_{1}-\cdots-x_{n-1}$ and $\mathbb{Z}_{+}$is the set of non-negative integers. The polynomials of degree less than or equal to $N$ are linearly spanned by $x^{v},|v|=\sum_{i=0}^{n-1} v_{i}=N$, which are linearly independent.

We denote by $\mathbb{Z}^{+}\left[x_{1}, \ldots, x_{n-1}\right]$ the cone generated by $x^{v}, v \in \mathbb{Z}_{+}^{n}$. By using this as a positive cone we define an order on $\mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$, i.e. $p_{1} \geq p_{2}$ if $p_{1}-p_{2} \in$ $\mathbb{Z}^{+}\left[x_{1}, \ldots, x_{n-1}\right]$. We call an element of $\mathbb{Z}^{+}\left[x_{1}, \ldots, x_{n-1}\right]$ positive.

We define

$$
\Lambda_{n}=\left\{\lambda=\left(\lambda_{1}, \ldots, \lambda_{n-1}\right) \in[0,1]^{n-1} \mid \sum_{j=1}^{n-1} \lambda_{j} \leq 1\right\}
$$

If $\lambda \in \Lambda_{n}$, then $\lambda_{0}$ denotes $1-\lambda_{1}-\cdots-\lambda_{n-1}$.
If $p \geq 0$, then $p$ is non-negative on $\Lambda_{n}$ and moreover satisfies that if $p(\lambda)=0$ for some $\lambda \in \Lambda_{n}$ then $p=0$ on the face of $\Lambda_{n}$ generated by $\lambda$. (This is because each $x^{v}$ satisfies this condition.) However, the converse does not follow if $n>2$ (compare this with the case $n=2$ in [24]).

For example, if $p \in \mathbb{Z}\left[x_{1}, x_{2}\right]$ is given by

$$
p\left(x_{1}, x_{2}\right)=x_{1}\left(x_{1}^{2}+\left(1-x_{1}-2 x_{2}\right)^{2}\right)
$$

then $p=0$ on the face $\Lambda^{\prime}=\left\{x \in \Lambda_{3} \mid x_{1}=0\right\}$, but $p>0$ on $\Lambda_{3} \backslash \Lambda^{\prime}$. If $p$ is expressed as

$$
p=\sum_{k+\ell+m=N} a_{k, \ell, m} x_{0}^{k} x_{1}^{\ell} x_{2}^{m},
$$

for some $N \in \mathbb{N}$ with $a_{k, \ell, m} \in \mathbb{Z}_{+}$, then that $\ell=0$ implies that $a_{k, \ell, m}=0$. Hence, we should have that

$$
x_{1}^{2}+\left(1-x_{1}-2 x_{2}\right)^{2}=\sum_{k+\ell+m=N-1} a_{k, \ell+1, m} x_{0}^{k} x_{1}^{\ell} x_{2}^{m}
$$

If $x_{1}=0$, the left-hand side vanishes for $x_{2}=\frac{1}{2}$ but not for $x_{2} \neq \frac{1}{2}$. This is not possible for the right-hand side.

Another example can be given by $p\left(x_{1}, x_{2}\right)=x_{1}^{2}+x_{2}^{2}-x_{1} x_{2}$. This is strictly positive on $\Lambda_{3}$ except for $x_{0}=1$ or $x_{1}=0=x_{2}$. If we express $p$ as $p=\sum_{k+\ell+m=N} a_{k, \ell, m} x_{0}^{k} x_{1}^{\ell} x_{2}^{m}$, then the highest order in $x_{0}$ must be $N-2$ and the sum of terms which include $x_{0}^{N-2}$ is $x_{0}^{N-2}\left(x_{1}^{2}+x_{2}^{2}-x_{1} x_{2}\right)=x_{0}^{N-2} p(x)$; thus $p$ is not positive.

Note that if $p \in \mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ is positive and $\Lambda^{\prime}$ is a face of $\Lambda_{n}$ with $\Lambda^{\prime} \cong \Lambda_{m}, 1<$ $m<n$, then $p \mid \Lambda^{\prime}$ can be understood as a positive element (or zero) of $\mathbb{Z}\left[x_{1}, \ldots, x_{m-1}\right]$ by taking $m-1$ free variables in a certain order.

The following is due to Handelman [9].
Proposition 2.1. The dimension group $K_{0}\left(A_{n}\right)$ is isomorphic to $\mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$, where the isomorphism is defined by

$$
\left[\prod_{m=-M}^{M} e_{i_{m} i_{m}}^{(m)}\right] \mapsto \prod_{m=-M}^{M} x_{i_{m}}
$$

for $i_{m}=0, \ldots, n-1$ and $M \in \mathbb{N}$, with $x_{0}=1-x_{1}-\cdots-x_{n-1}$.
The positive cone $C=\mathbb{Z}^{+}\left[x_{1}, \ldots, x_{n-1}\right]$ is the cone with the following properties:
(1) $x_{i} \in C$ for $i=0, \ldots, n-1$, where $x_{0}=1-x_{1}-\cdots-x_{n-1}$;
(2) $p q \in C$ if $p, q \in C$;
(3) $p \in C$ if $p \in \mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ is strictly positive as a function on $\Lambda_{n}$.

Proof. In this proof we write $A_{n}$ as $A$ omitting the subscript $n$.
Let $\left(Z_{k}\right)$ be an increasing sequence of finite subsets of $\mathbb{Z}$ such that the number $\left|Z_{k}\right|$ of elements in $Z_{k}$ is $k$ and $\bigcup_{k} Z_{k}=\mathbb{Z}$, and let $A_{0}=\mathbb{C} 1$ and

$$
A_{k}=\left(\bigotimes_{m \in Z_{k}} M_{n}^{(m)}\right)^{\gamma}
$$

Then $\left(A_{k}\right)$ is an increasing sequence of finite-dimensional $C^{*}$-subalgebras of $A$ such that $A=\overline{\bigcup_{k} A_{k}}$.

Let $e_{0}, \ldots, e_{n-1}$ be the canonical basis for $\mathbb{Z}^{n}$ and let, for each $k \in \mathbb{N}$,

$$
V_{k}=\left\{v \in \mathbb{Z}_{+}^{n}| | v \mid=\sum_{i=0}^{n-1} v_{i}=k\right\}
$$

where $\mathbb{Z}_{+}=\{0,1,2, \ldots$,$\} . For v \in V_{k}$ we denote by $N(v)$ the number of sequences $\left(e_{i_{1}}, \ldots, e_{i_{k}}\right)$ in $\left\{e_{0}, \ldots, e_{n-1}\right\}$ such that $v=\sum_{j=1}^{k} e_{i_{j}}$, i.e.

$$
N(v)=|v|!/ v!=|v|!\left(v_{0}!v_{1}!\cdots v_{n-1}!\right)^{-1}
$$

Define

$$
E_{v}=\sum \Pi_{m \in Z_{k}} e_{i_{m}, i_{m}}^{(m)} \in A_{k}
$$

where the sum is taken over all the map $i: Z_{k} \rightarrow\{0, \ldots, n-1\}$ such that $\sum_{m \in Z_{k}} e_{i_{m}}=v$. It follows that $E_{v}$ is a minimal central projection of $A_{k}$ and $A_{k} E_{v} \cong$ $M_{N(v)}$. The embedding of $A_{k}$ into $A_{k+1}$ is given as follows: for $v \in V_{k}$ and $w \in V_{k+1}$, $A_{k} E_{v}$ is mapped into $A_{k+1} E_{w}$ with multiplicity one if and only if $w=v+e_{i}$ for some $i$ or $w \geq v$ in the sense that $w_{i} \geq v_{i}$ for all $i$.

We define a map $\psi_{k}$ of $K_{0}\left(A_{k}\right)$ into $\mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ by the following. If $e$ is a minimal projection in $A_{k} E_{v}$, then

$$
\psi_{k}([e])=x^{v}=x_{0}^{v_{0}} x_{1}^{v_{1}} \cdots x_{n-1}^{v_{n-1}}
$$

where $v=\left(v_{0}, \ldots, v_{n-1}\right)$ and $x_{0}=1-x_{1}-\cdots-x_{n-1}$. Since $\psi_{k+1} \circ \iota_{*}$ coincides with $\psi_{k}$ on $K_{0}\left(A_{k}\right)$, where $\iota$ is the embedding of $A_{k}$ into $A_{k+1}$, we can define a homomorphism $\psi$ of $K_{0}(A)$ into $\mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$. Since the range of $\psi_{k}$ is the polynomials of order less than or equal to $k$, it follows that $\psi$ is surjective. Since $\psi_{k}$ is injective, it also follows that $\psi$ is injective. It is easy to check that $\psi$ is defined as indicated in the statement and that the range of $\psi$ on the projections in $A$ generate the positive cone $\mathbb{Z}^{+}\left[x_{1}, \ldots, x_{n-1}\right]$. Therefore, $\psi$ is indeed the required isomorphism. Note that $\psi([1])=1$.

We have to show the statement on the positive cone. It is immediate that conditions (1) and (2) are valid and are enough to generate $\mathbb{Z}^{+}\left[x_{1}, \ldots, x_{n-1}\right]$. We now show the validity of condition (3), which has been known for some time at least for $n=2$ (see, e.g., [11, p. 126]).

First of all note that for each $\lambda \in \Lambda_{n}$ there is a unique tracial state $\tau_{\lambda}$ on $A$ such that for a minimal projection $e \in A_{k} E_{v}$ with $v \in V_{k}$,

$$
\tau_{\lambda}(e)=\lambda^{v}=\lambda_{0}^{v_{0}} \lambda_{1}^{v_{1}} \cdots \lambda_{n-1}^{v_{n-1}} .
$$

Moreover, $\tau_{\lambda}$ is factorial and all the factorial tracial states of $A$ are of this form, as shown by the following lemma.

Let $p \in \mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ be such that $p$ is strictly positive on $\Lambda_{n}$. Since $\psi$ is a group isomorphism, there are projections $e, f$ in $M_{m} \otimes A_{k}$ for some $m, k$ such that $p=\psi([e]-[f])$. What we have to show is that $[e] \geq[f]$, i.e. $f$ is equivalent to a subprojection of $e$.

Suppose the contrary, i.e. there is a sequence $\left(v_{\ell}\right)_{\ell>k}$ with $v_{\ell} \in V_{\ell}$ such that

$$
\operatorname{rank}\left(e\left(1_{m} \otimes E_{v_{\ell}}\right)\right)<\operatorname{rank}\left(f\left(1 \otimes E_{v_{\ell}}\right)\right)
$$

in $M_{m} \otimes A_{\ell} E_{v_{\ell}}$. We extend the state on $A_{\ell}$ defined by

$$
x \mapsto N\left(v_{\ell}\right)^{-1} \operatorname{Tr}\left(x E_{v_{\ell}}\right)
$$

to a state $\varphi_{\ell}$ on $A$, where $\operatorname{Tr}$ is the trace on $A_{\ell} E_{v_{\ell}} \cong M_{N\left(v_{\ell}\right)}$. We take a weak* limit point $\tau$ of $\left(\varphi_{\ell}\right)_{\ell>k}$. Then it follows that $\tau$ is a tracial state on $A$ for which

$$
\operatorname{Tr}_{m} \otimes \tau(e) \leq \operatorname{Tr}_{m} \otimes \tau(f),
$$

where $\operatorname{Tr}_{m}$ is the trace on $M_{m}$. Since $\tau$ belongs to the closed convex hull of $\tau_{\lambda}, \lambda \in \Lambda_{n}$, this contradicts the assumption that $\operatorname{Tr}_{m} \otimes \tau_{\lambda}(e)>\operatorname{Tr}_{m} \otimes \tau_{\lambda}(f)$ for all $\lambda \in \Lambda_{n}$. This concludes the proof of condition (3). (We see in Lemma 2.3 that the above $\tau$ is actually $\tau_{\lambda}$ for some $\lambda \in \Lambda_{n}$.)

Lemma 2.2. Let $\tau$ be a factorial tracial state on $A_{n}$. Then there is $a \lambda \in \Lambda_{n}$ such that $\tau=\tau_{\lambda}$.

Proof. Since $\tau$ is a tracial state on $A=A_{n}$, there is map $\varphi_{k}: V_{k} \rightarrow[0,1]$ such that $\tau(e)=\varphi_{k}(v)$ for a minimal projection $e$ in $A_{k} E_{v}$. We should note that $\left(\varphi_{k}\right)$ is independent of the choice of $\left(Z_{k}\right)$ by which $A_{k}$ is defined as $A_{k}=\left(\bigotimes_{m \in Z_{k}} M_{n}^{(m)}\right)^{\gamma}$. By choosing $Z_{k+1}=Z_{k} \cup\{N\}$ for a large $N$ and by using

$$
\varphi_{k+1}\left(v+e_{i}\right)=\tau\left(e e_{i i}^{(N)}\right)
$$

with $e$ a minimal projection in $A_{k} E_{v}$, we get that

$$
\varphi_{k+1}\left(v+e_{i}\right)=\varphi_{k}(v) \varphi_{1}\left(e_{i}\right)
$$

for $v \in V_{k}$. Here we have used the assumption that $\tau$ is factorial and that $\left(e_{i i}^{(N)}\right)_{N}$ is a central sequence with $\tau\left(e_{i i}^{(N)}\right)=\varphi_{1}\left(e_{i}\right)$. Setting $\lambda_{i}=\varphi_{1}\left(e_{i}\right)$ for $i=0, \ldots, n-1$, we get that $\sum_{i=0}^{n-1} \lambda_{i}=1$ and $\varphi_{k}(v)=\lambda^{v}=\lambda_{0}^{v_{0}} \cdots \lambda_{n-1}^{v_{n-1}}$. Hence, $\lambda=\left(\lambda_{1}, \ldots, \lambda_{n-1}\right) \in \Lambda_{n}$ and $\tau=\tau_{\lambda}$.

In the following lemma we adopt the notation in the proof of the above proposition. In particular, we simply denote by $A$ the $C^{*}$-algebra $A_{n}$ and we express $A$ as the closure of $\bigcup_{k} A_{k}$, where $A_{k}=\left(\otimes_{Z_{k}} M_{n}\right)^{\gamma}$ and $\left|Z_{k}\right|=k$.
LEMMA 2.3. Let $\left(m_{\ell}\right)$ be an increasing sequence of integers, $\left(v_{\ell}\right)$ a sequence in $\bigcup_{m} V_{m}$, and $\lambda \in \Lambda_{n}$ such that $v_{\ell} \in V_{m_{\ell}}, m_{\ell} \rightarrow \infty$, and

$$
\frac{v_{\ell, i}}{m_{\ell}} \rightarrow \lambda_{i}
$$

as $\ell \rightarrow \infty$ for $i=0,1, \ldots, n-1$, where $v_{\ell}=\sum_{i=0}^{n-1} v_{\ell, i} e_{i}$. Let $\varphi_{\ell}$ be a state of $A$ such that $\varphi_{\ell} \mid E_{v_{\ell}} A_{m_{\ell}}=N\left(v_{\ell}\right)^{-1} \mathrm{Tr}$, where $\operatorname{Tr}$ is the trace on $M_{N\left(v_{\ell}\right)} \cong E_{v_{\ell}} A_{m_{\ell}}$. Then $\varphi_{\ell}$ converges to $\tau_{\lambda}$ as $\ell \rightarrow \infty$.

Proof. Let $k \in \mathbb{N}$. We evaluate $\varphi_{\ell} \mid A_{k}$ for $m_{\ell} \gg k$. Note that $\varphi_{\ell} \mid A_{k}$ is a tracial state; hence, it is expressed as

$$
\varphi_{\ell}\left|A_{k}=\sum_{w \in V_{k}} d_{w} N(w)^{-1} \operatorname{Tr}\right| E_{w} A_{k}
$$

for non-negative constants $d_{w}$ with $\sum_{w} d_{w}=1$.
The multiplicity with which $E_{w} A_{k}$ is embedded into $E_{v_{\ell}} A_{m_{\ell}}$ is

$$
N\left(w, v_{\ell}\right)=\frac{\left(m_{\ell}-k\right)!}{\left(v_{\ell, 0}-w_{0}\right)!\left(v_{\ell, 1}-w_{1}\right)!\cdots\left(v_{\ell, n-1}-w_{n-1}\right)!}
$$

if $w \leq v_{\ell}$; otherwise it is zero. (For $w \leq v_{\ell}, N\left(w, v_{\ell}\right)$ is the coefficient of $x^{v_{\ell}-w}$ in the expansion of $\left(x_{0}+x_{1}+\cdots+x_{n-1}\right)^{m_{\ell}-\bar{k}}$.) Hence, we have that if $d_{w}$ is non-zero, then

$$
d_{w}=\varphi_{\ell}\left(E_{w}\right)=N\left(v_{\ell}\right)^{-1} N(w) N\left(w, v_{\ell}\right)
$$

Suppose that all $\lambda_{i}$ are positive. If $m_{\ell}$ is so large that $v_{\ell, i} \gg w_{i}$, then it follows that

$$
d_{w}=N(w) \frac{\left(m_{\ell}-k\right)!v_{\ell, 0}!\cdots v_{\ell, n-1}!}{m_{\ell}!\left(v_{\ell, 0}-w_{0}\right)!\cdots\left(v_{\ell, n-1}-w_{n-1}\right)!}
$$

is approximately equal to

$$
N(w) \lambda_{0}^{w_{0}} \lambda_{1}^{w_{1}} \cdots \lambda_{n-1}^{w_{n-1}} .
$$

This shows that $\varphi_{\ell} \mid A_{k}$ converges to $\tau_{\lambda} \mid A_{k}$.
Next we consider the case where some of $\lambda_{i}$ are zero. For simplicity suppose that $\lambda_{i}=0$ for $i=j, \ldots, n-1$ with some $j>0$. Then we have that for $w \in V_{k}$ with $w_{i}=0$, $i=j, \ldots, n-1$,

$$
d_{w} \approx N(w) \lambda_{0}^{w_{0}} \lambda_{1}^{w_{1}} \cdots \lambda_{j-1}^{w_{j-1}},
$$

and for $w \in V_{k}$ with $\sum_{i=j}^{n-1} w_{i}=s>0, d_{w}$ is of the order of $m_{\ell}^{-s} v_{j}^{w_{j}} \cdots v_{n-1}^{w_{n-1}}$, which converges to zero. Thus, in this case too one can conclude that $\varphi_{\ell} \mid A_{k}$ converges to $\tau_{\lambda} \mid A_{k}$. Since $k$ is arbitrary, it follows that $\left(\varphi_{\ell}\right)$ converges to $\tau_{\lambda}$ in the weak* topology.

Suppose that $p \in \mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ is strictly positive as a function on $\Lambda_{n}$ except for some extreme points. If the degree of $p$ is less than or equal to $k$, then $p$ can be expressed as

$$
p\left(x_{1}, \ldots, x_{n-1}\right)=\sum_{w \in V_{k}} c_{w} x^{w}
$$

in a unique way, where the coefficients $c_{w}$ are integers and $x^{w}=x_{0}^{w_{0}} x_{1}^{w_{1}} \cdots x_{n-1}^{w_{n-1}}$ with $x_{0}=1-x_{1}-\cdots-x_{n-1}$. To obtain a similar expression in terms of $x^{v}, v \in V_{k+1}$, we just have to multiply the right-hand side with $x_{0}+x_{1}+\cdots+x_{n-1}$ and expand it.

Let $a_{0}=\max \left\{w_{0} \mid c_{w} \neq 0\right\}$ and define

$$
q_{0}\left(x_{1}, \ldots, x_{n-1}\right)=\sum_{w \in V_{k}, w_{0}=a_{0}} c_{w} x_{1}^{w_{1}} \cdots x_{n-1}^{w_{n-1}}
$$

We note that $q_{0}$ does not depend on $k$ (if $k$ increases by one, then so does $a_{0}$ ); unique to the vertex $x_{0}=1$. If $p(0,0, \ldots, 0)=c_{(k, 0, \ldots, 0)}=c \neq 0$, then $q_{0}=c$. Similarly we can define $q_{i}$ as a polynomial in $x_{j}, j \neq i$. If $p$ is positive in $\mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$, it follows that $q_{i}, i=0, \ldots, n-1$ are positive (or zero) in $\mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$.

More generally we define the 0 -vanguard $V_{k}^{0}$ of $p$ on the level $k$ as the set of $w \in V_{k}$ which satisfies that $c_{w} \neq 0$ and that if $v$ is in front of $w$ in the direction to the zeroth vertex, i.e. $v=w+s e_{0}-e_{i_{1}}-e_{i_{2}}-\cdots-e_{i_{s}} \in V_{k}$ with $i_{j} \neq 0$ for some $s>0$, then $c_{v}=0$.

Then the map $w \mapsto w+e_{0}$ from $V_{k}$ into $V_{k+1}$ restricts to a bijection from $V_{k}^{0}$ onto $V_{k+1}^{0}$ preserving the coefficients $c_{w}$.

To prove this, let $w \in V_{k}^{0}$. If $v^{\prime} \in V_{k+1}$ is in front of $w+e_{0}$, then $v^{\prime}-e_{0} \in V_{k}$ is in front of $w$ and hence $c_{v^{\prime}-e_{0}}=0$. If $v^{\prime}=w^{\prime}+e_{i}$ for some $i \neq 0$, then $w^{\prime}=v^{\prime}-e_{i}$ is in front of $w+e_{0}-e_{i}$, which is in front of $w$, and hence $c_{w^{\prime}}=0$. This way we can conclude that $c_{v^{\prime}}=0$ as $c_{v^{\prime}}$ is the sum of $c_{v^{\prime}-e_{i}}$ with $v^{\prime}-e_{i} \in V_{k}$. If $w+e_{0}=w^{\prime}+e_{i}$ for some $w^{\prime} \neq w$, then $w^{\prime}=w+e_{0}-e_{i}$ is in front of $w$ and hence $c_{w^{\prime}}=0$. This implies that $c_{w+e_{0}}=c_{w} \neq 0$. Hence, it follows that $w+e_{0} \in V_{k+1}^{0}$ and $c_{w+e_{0}}=c_{w}$.

On the other hand let $v \in V_{k+1}^{0}$; if $w=v-e_{0} \notin V_{k}^{0}$, then it follows that $c_{w}=0$ or there is a $w^{\prime}(\neq w)$ in front of $w$ such that $c_{w^{\prime}} \neq 0$. Since $c_{v} \neq 0$, there is $i$ such that $c_{v-e_{i}} \neq 0$. If $i=0$, then it means that $c_{w} \neq 0$. If $i \neq 0$, then $v-e_{i}=w+e_{0}-e_{i}$ is in front of $w$ and $c_{w+e_{0}-e_{i}} \neq 0$. Thus, in either case there is a $w^{\prime}$ in front of $w$ such that $c_{w^{\prime}} \neq 0$. We choose such a $w^{\prime}$ of maximal $w_{0}^{\prime}$, i.e. $w^{\prime} \in V_{k}^{0}$. Then $w^{\prime}+e_{0}$ is in front
of $w+e_{0}=v$, which implies that $c_{w^{\prime}+e_{0}}=0$. If $w^{\prime}+e_{0}=u+e_{i}$ for some $i \neq 0$, $u=w^{\prime}+e_{0}-e_{i}$ is in front of $w^{\prime}$ and hence $c_{u}=0$. Since $c_{w^{\prime}+e_{0}}=c_{w^{\prime}}+\sum_{i=1}^{n-1} c_{w^{\prime}+e_{0}-e_{i}}$, this would imply that $c_{w^{\prime}+e_{0}}=c_{w^{\prime}}$, which is a contradiction. This proves that $v-e_{0} \in V_{k}^{0}$ for $v \in V_{k+1}^{0}$. From the first part it also follows that $c_{v-e_{0}}=c_{v}$. Hence, the above assertion is now shown.

One can similarly define the $i$-vanguard $V_{k}^{i}$ of $p$ on the level $k$ for $i=1, \ldots, n-1$ and we have a similar bijection from $V_{k}^{i}$ onto $V_{k+1}^{i}$ with the property $c_{w+e_{i}}=c_{w}, w \in V_{k}^{i}$.

Definition 2.4. Let $p \in \mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ and $i=0,1, \ldots, n-1$. The $i$-vanguard $V^{i}(p)$ of $p$ is the set of $v:\{0,1, \ldots, n-1\} \backslash\{i\} \rightarrow \mathbb{Z}$ such that $v^{(k)} \in V_{k}^{i}(p)$, where $V_{k}^{i}(p)$ is the $i$ th vanguard of $p$ on the level $k$ and $v^{(k)}$ is defined by $v_{j}^{(k)}=v_{j}$ for $j \neq i$ and $v_{i}^{(k)}=k-\sum_{j} v_{j}$. When $p=\sum_{v \in V_{k}} c_{v} x^{v}$ and $v \in V^{i}(p)$ for some $i$, we set $p(v)=c_{v}$.

The following gives a sufficient condition for positivity.
Proposition 2.5. Let $p \in \mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ be such that $p$ is strictly positive on $\Lambda_{n}$ except for some extreme points. Then $p$ is positive if and only if $p(v)>0$ for all $v \in V^{i}(p), i=0,1, \ldots, n-1$.

Proof. We write $p$ as

$$
p\left(x_{1}, \ldots, x_{n-1}\right)=\sum_{w \in V_{k}} c_{w} x^{w}
$$

where $k$ is greater than or equal to the degree of $p$. We have to show that for a sufficiently large $k$, all $c_{w}$ are non-negative by using the assumption that $p$ is strictly positive on $\Lambda_{n}$ except for the vertices and that $c_{w}>0$ for $w \in \bigcup_{i=0}^{n-1} V_{k}^{i}(p)$; the other implication is obvious.

If the assertion is false, there is a sequence $\left(v_{\ell}\right)$ such that $v_{\ell} \in V_{\ell}$ and $c_{v_{\ell}}<0$. We may suppose that $v_{\ell, i} / \ell \rightarrow \lambda_{i}$ for each $i=0,1, \ldots, n-1$. Then the above lemma shows that

$$
c_{v_{\ell}} N\left(v_{\ell}\right)^{-1} \rightarrow p(\lambda) .
$$

where $N\left(v_{\ell}\right)=\ell!\left(v_{\ell}!\right)^{-1}$ and $\lambda=\left(\lambda_{1}, \ldots, \lambda_{n-1}\right) \in \Lambda_{n}$. (To make the correspondence with the lemma clearer, let us find two projections $e, f$ in (a matrix-tensored) $A$ such that $p=[e]-[f]$ and let $\varphi_{\ell}$ be a state of $A$ such that $\varphi_{\ell} \mid E_{v_{\ell}} A_{\ell}$ is a tracial state. Then it is shown that $\varphi_{\ell}(e-f)=c_{v_{\ell}} N\left(v_{\ell}\right)^{-1}$, which is negative, converges to $\tau_{\lambda}(e-f)=p(\lambda)$.) If $\lambda$ is not a vertex of $\Lambda_{n}$, then $p(\lambda)>0$, which implies that $\lambda$ must be a vertex; we may thus assume that $v_{\ell, 0} / \ell \rightarrow 1$ and $v_{\ell, i} / \ell \rightarrow 0$ for $i=1, \ldots, n-1$.

Let $w \in V_{k}$. If $w \leq v_{\ell}, E_{v} A_{k}$ is embedded into $E_{v_{\ell}} A_{v_{\ell}}$ with the multiplicity $N\left(w, v_{\ell}\right)$ defined before. Hence, we have that

$$
c_{v_{\ell}}=\sum_{w \in V_{k}, w \leq v_{\ell}} c_{w} N\left(w, v_{\ell}\right)
$$

If $c_{w}<0$ for some $w \in V_{k}$ with $w \leq v_{\ell}$, then there is a $w^{\prime} \in V_{k}^{0}(p)$ such that $w^{\prime}$ is in front of $w$ (and $c_{w}>0$ ). Then we can argue that $c_{v_{\ell}}$ must be positive for large $\ell$ by showing that the contribution of $c_{w} N\left(w, v_{\ell}\right)$ to $c_{v_{\ell}}$, which is negative, is overshadowed by a tiny
portion of $c_{w^{\prime}} N\left(w^{\prime}, v_{\ell}\right)$ which is positive, as follows. Since

$$
N\left(w, v_{\ell}\right)=\frac{(\ell-k)!}{\left(v_{\ell, 0}-w_{0}\right)!\cdots\left(v_{\ell, n-1}-w_{n-1}\right)!},
$$

and $w_{0}^{\prime}>w_{0}$ and $w_{i}^{\prime} \leq w_{i}$ for $i>0$, we have that $w^{\prime} \leq v_{\ell}$ for large $\ell$ and

$$
\frac{N\left(w^{\prime}, v_{\ell}\right)}{N\left(w, v_{\ell}\right)}=\frac{\left(v_{\ell, 0}-w_{0}\right)!}{\left(v_{\ell, 0}-w_{0}^{\prime}\right)!} \frac{\left(v_{\ell, 1}-w_{1}\right)!\cdots\left(v_{\ell, n-1}-w_{n-1}\right)!}{\left(v_{\ell, 1}-w_{1}^{\prime}\right)!\cdots\left(v_{\ell, n-1}-w_{n-1}^{\prime}\right)!}
$$

is approximately equal to or more than

$$
\frac{\ell^{w_{0}^{\prime}-w_{0}}}{v_{\ell, 1}^{w_{1}-w_{1}^{\prime}} \cdots v_{\ell, n-1}^{w_{n-1}-w_{n-1}^{\prime}}}
$$

which tends to infinity as $\ell \rightarrow \infty$. Hence, we can conclude that $c_{v_{\ell}}>0$ for a sufficiently large $\ell$, which is a contradiction.

LEMMA 2.6. Let $p$ and $q$ be positive elements of $\mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ and $i=0,1, \ldots, n-1$. Then the $i$-vanguard $V^{i}(p q)$ of $p q$ is contained in $\left\{v+w \mid v \in V^{i}(p), w \in V^{i}(q)\right\}$.
Proof. Let $p=\sum_{v \in V_{K}} c_{v} x^{v}$ and $q=\sum_{v \in V_{K}} d_{v} x^{v}$ for some $K \in \mathbb{N}$. Suppose that there are $v, w \in V_{K}$ such that $c_{v}>0, d_{w}>0$, and $v+w$ is in the $i$-vanguard of $p q$ on the level $2 K$. If $v \notin V_{K}^{i}(p)$, then there is a $v^{\prime} \in V_{K}$ such that $v^{\prime}$ is in front of $v$ toward the $i$ th vertex and $c_{v^{\prime}} \neq 0$. Then $v^{\prime}+w$ is in front of $v+w$ toward the $i$ th vertex and the coefficient of $x^{v^{\prime}+w}$ in the expansion of $p q$ in $x^{\mu}, \mu \in V_{2 K}$ is greater than or equal to $c_{v^{\prime}} d_{w}$, i.e. does not vanish, which is a contradiction. Similarly we can conclude that $w \in V_{K}^{i}(q)$. If $u \in V_{2 K}$ satisfies that if $u=v+w$ with $v, w \in V_{K}$ then $c_{v}=0$ or $d_{w}=0$, then $u$ cannot belong to the vanguards of $p q$. This concludes the proof.

Later we use a more elaborate form of the following.
COROLLARY 2.7. Let $p \in \mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ be a positive element such that $p$ is strictly positive on $\Lambda_{n}$ except for the vertices. Then there exists a $K \in \mathbb{N}$ such that $p-q$ is positive for any $q \in\left(x_{0} x_{1} \cdots x_{n-1}\right)^{K} \mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ if $p-q$ is strictly positive on $\Lambda_{n}$ except for the vertices.

Proof. Suppose that the degree of $p$ is less than $K$. Let $q \in\left(x_{0} \cdots x_{n-1}\right)^{K} \mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$. We express $p$ and $q$ as

$$
p=\sum_{v \in V_{L}} c_{v} x^{v}, \quad q=\sum_{v \in V_{L}} d_{v} x^{v}
$$

for some $L \geq K$. We consider the $i$-vanguard $V_{L}^{i}(p-q)$ of $p-q$ on the level $L$. Note that the $i$-vanguard $V_{L}^{i}(q)$ is confined to $\left\{v \mid v_{j} \geq K\right.$ for $\left.j \neq i\right\}$.

Let $v \in V_{L}^{i}(p-q)$ and suppose that $v_{i} \leq L-K$. Since the $i$-vanguard of $p$ on the level $L$ is $V_{K-1}^{i}(p)+(L-K+1) e_{i}$, we have that $d_{v} \neq 0$. (If $d_{v}=0$, then we must have that $c_{v} \neq 0$; since $v$ does not belong to $V_{L}^{i}(p)$, there is $w \in V_{L}^{i}(p)$ in front of $v$ toward the $i$ th vertex such that $c_{w} \neq 0$ and $w_{i}>L-K$. Since $d_{w}=0$ for such a $w$, we have that $c_{w}-d_{w}=c_{w} \neq 0$, which is a contradiction.) Since $d_{v} \neq 0$, we have that all $v_{j} \geq K$. If $w \in V_{L}$ satisfies that $w_{i}=L-K+1$, then $w_{i}-v_{i}>0$ and $w_{j}-v_{j}<0$
for $j \neq i$, i.e. we have that all $w \in V_{L}$ with $w_{i}=L-K+1$ are in front of $v$ towards the $i$ th vertex. Thus, all of the points $w \in V_{L}$ with $w_{i}>L-K$ are in front of $v$; hence, $v$ cannot belong to $V_{L}^{i}(p-q)$. This implies that if $v \in V_{L}^{i}(p-q)$, then $v_{i}>L-K$; hence, $V_{L}^{i}(p-q)=V_{K-1}^{i}(p)+L-K+1$. Note that $c_{w}-d_{w}=c_{w}>0$ for $w \in V_{L}^{i}(p-q)$. Since $p-q$ is strictly positive on $\Lambda_{n} \backslash \operatorname{Ex}\left(\Lambda_{n}\right)$, the conclusion follows from the previous proposition.

To estimate the vanguards we use the following lemma.
Lemma 2.8. Let $p$ be a non-negative element of $\mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ such that

$$
p\left(x_{1}, \ldots, x_{n-1}\right) \geq c \sum_{i \neq j}\left(x_{i} x_{j}\right)^{S}
$$

on $\Lambda_{n}$ for some $c>0$ and $S \in \mathbb{N}$. Then the vanguards $V^{i}(p)$ of $p$ are confined in the $S$-neighborhood of vertices; more precisely, if $\omega \in V^{i}(p)$ for some $i=0,1, \ldots, n-1$, then there is a $j \neq i$ such that $\omega_{k}=0$ for $k \neq i, j$ and $\omega_{j} \leq S$.

Proof. Let $p=\sum_{v \in V_{K}} c_{v} x^{v}$ and let $i=0,1, \ldots, n-1$ and $j \neq i$. If $x \in \Lambda_{n}$ satisfies that $x_{k}=0$ for $k \neq i, j$, i.e. if $x$ is on the edge $[i, j]$ between the vertices $x_{i}=1$ and $x_{j}=1$, then $p(x) \geq c\left(x_{i} x_{j}\right)^{S}$. Since for $x_{i} \approx 1\left(\right.$ and $\left.x_{j}=1-x_{i} \approx 0\right)$,

$$
p(x)=\sum_{w \in V_{k}^{i}(p)} c_{w} x^{w}+\cdots
$$

there must be an $\omega \in V_{k}^{i}(p)$ such that $\omega_{k}=0$ for $k \neq i, j$ and $\omega_{j} \leq S$. Hence, $V_{K}^{i}(p)$ contains $\omega^{(j)}$ for each $j \neq i$ such that $\omega_{j}^{(j)} \leq S$ and $\omega_{k}^{(j)}=0$ for $k \neq i, j$ (and $\omega_{i}^{(j)}=K-\omega_{j}^{(j)}$ ). This implies that if $\omega \in V_{K}$ has $j \neq i$ such that $\omega_{j} \geq S$, then $\omega^{(j)}$ is in front of $\omega$ (towards the $i$ th vertex), i.e. $\omega \notin V_{K}^{i}(p)$ unless $\omega=\omega^{(j)}$. This proves the assertion.

Lemma 2.9. Let $k=\left(k_{0}, k_{1}, \ldots, k_{n-1}\right) \in \mathbb{Z}_{+}^{n}$. Then the order ideal of $\mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ generated by $x^{k}=x_{0}^{k_{0}} \cdots x_{n-1}^{k_{n-1}}$ is equal to $x^{k} \mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$, where $x_{0}=1-x_{1}-\cdots-$ $x_{n-1}$.

Proof. Let $p$ be a non-negative element of $\mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ such that $N x^{k}-p \geq 0$ for some $N \in \mathbb{N}$. Then for a sufficiently large $L \in \mathbb{N}$, we have $c, d: V_{L} \rightarrow \mathbb{Z}_{+}$such that

$$
N x^{k}=\sum_{w \in V_{L}} c_{w} x^{w}, \quad p=\sum_{w \in V_{L}} d_{w} x^{w}
$$

and $c_{w} \geq d_{w} \geq 0$ for all $w \in V_{L}$. Since $c_{w}>0$ if and only if $w \geq k$, this implies that if $d_{w}>0$ then $w \geq k$ and that

$$
p=x^{k}\left(\sum_{w \in V_{L}, w \geq k} d_{w} x^{w-k}\right) .
$$

Since the other inclusion is obvious, this concludes the proof.

For $k \in \mathbb{Z}_{+}^{n}$ we denote by $I(k)$ the ideal of $A_{n}$ corresponding to the order ideal of $K_{0}\left(A_{n}\right) \cong \mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ generated by $x^{k}=x_{0}^{k_{0}} x_{1}^{k_{1}} \cdots x_{n-1}^{k_{n-1}}$. The ideal $I_{i}$ which is used to define the character $\varphi_{i}$ in the main Theorem 1.1 is given as

$$
I_{i}=\sum_{j \neq i} I\left(e_{j}\right) .
$$

We note that $I(k) \cap I(\ell)=I(k \vee \ell)$ for $k, \ell \in \mathbb{Z}_{+}^{n}$, where $(k \vee \ell)_{i}=\max \left(k_{i}, \ell_{i}\right)$.

## 3. Covariant irreducible representations

In this section we define some covariant irreducible representations of $\left(A_{n}, \sigma\right)$, which are an extension of what is well known for the case $n=2$ and will be used in the proof of Lemma 4.4, the very first step of the induction.

Let $k=\left(k_{1}, \ldots, k_{n-1}\right) \in \mathbb{Z}_{+}^{n-1}$ such that $|k|=\sum_{i} k_{i} \neq 0$ and let

$$
\Omega_{k}=\left\{\left(S_{1}, \ldots, S_{n-1}\right) \in \mathcal{P}(\mathbb{Z})^{n-1}| | S_{i} \mid=k_{i}, i \neq j \Rightarrow S_{i} \cap S_{j}=\emptyset\right\} .
$$

For $S, T \in \Omega_{k}$ let

$$
e(S, T)=\prod_{i=1}^{n-1}\left(\prod_{m \in S_{i}} e_{i 0}^{(m)} \prod_{m \in T_{i}} e_{0 i}^{(m)}\right) \in A_{n} .
$$

Note that $e(S, T) \in I\left(0, k_{1}, \ldots, k_{n-1}\right)=I(0, k)$. On the Hilbert space $\ell^{2}\left(\Omega_{k}\right)$ we define a representation $\pi_{k}$ of $I(0, k)$ such that $\operatorname{Ker} \pi_{k} \supset I\left((0, k)+e_{i}\right)$ for $i=1,2, \ldots, n-1$ and

$$
\pi_{k}(e(S, T)) \xi_{U}=\delta_{T, U} \xi_{S}
$$

for $S, T, U \in \Omega_{k}$, where $\left(\xi_{U}\right)_{U \in \Omega_{k}}$ is the canonical basis of $\ell^{2}\left(\Omega_{k}\right)$. (In the notation of the proof of Proposition 2.1, $I(0, k) \cap A_{\ell}$ is the ideal of $A_{\ell}=\left(\otimes_{Z_{\ell}} M_{n}\right)^{\gamma}$ generated by $E_{v}, v \in V_{\ell}$ with $v \geq(0, k)$ and $\pi_{k}(x)=\pi_{k}\left(x E_{\hat{k}}\right), x \in A_{\ell}$, where $\hat{k}=(\ell-|k|, k) \in V_{\ell}$. For $x \in A_{\ell}$ it follows that there is $c(S, T) \in \mathbb{C}$ for each pair $S, T \in \Omega_{k}$ with $\cup S, \cup T \subset Z_{\ell}$ such that $e(S, S) x E_{\hat{k}} e(T, T)=c(S, T) e(S, T) E_{\hat{k}}$ and $\sum e(S, S) x E_{\hat{k}} e(T, T)=x E_{\hat{k}}$, where the sum is taken over all pairs $S, T$ with $\cup S, \cup T \subset Z_{\ell}$.) Note that $\pi_{k}(I(0, k))$ is $\mathcal{K}=\mathcal{K}\left(\ell^{2}\left(\Omega_{k}\right)\right)$, the $C^{*}$-algebra of compact operators on $\ell^{2}\left(\Omega_{k}\right)$ and that $\pi_{k}(I(m, k))=$ $\mathcal{K}\left(\ell^{2}\left(\Omega_{k}\right)\right)$ for any $m \in \mathbb{N}$ (because $I(m, k)$ is an ideal of $I(0, k)$ not contained in the kernel of $\pi_{k}$ ). Since $I(0, k)$ is an ideal of $A_{n}, \pi_{k}$ naturally extends to an irreducible representation of $A_{n}$, which again is denoted by $\pi_{k}$. For $k=0=(0, \ldots, 0)$, we let $\pi_{0}$ denote the character $\varphi_{0}$ on $A_{n}$.

Lemma 3.1. For $k \in \mathbb{Z}_{+}^{n-1}$ define an irreducible representation $\pi_{k}$ of $A_{n}$ as above. Then

$$
\operatorname{Ker} \pi_{k}=\sum_{i=1}^{n-1} I\left(\left(k_{i}+1\right) e_{i}\right) .
$$

Proof. If $k=0$, then $\pi_{0}=\varphi_{0}$ is a character. We have already noted this at the end of the previous section.

If $k \neq 0$, then $\pi_{k}(I(0, k))=\mathcal{K}$ and $\operatorname{Ker}\left(\pi_{k} \mid I(0, k)\right)=\sum_{i=1}^{n-1} K_{0}(I(0, k)) x_{i}$, because $K_{0}(I(0, k))$ includes the latter as a maximal ideal. Hence, Ker $\pi_{k}$ is the largest ideal
$J$ of $A_{n}$ such that $K_{0}(J) \cap K_{0}(I(0, k))=\sum_{i=1}^{n-1} K_{0}(I(0, k)) x_{i}$. It is obvious that $J \supset \sum_{i=1}^{n-1} I\left(\left(k_{i}+1\right) e_{i}\right)$.

If $J \not \subset \sum_{i=1}^{n-1} I\left(\left(k_{i}+1\right) e_{i}\right)$, then there is a positive $p \in K_{0}(J)$ such that

$$
p \notin \sum_{i=1}^{n-1} K_{0}\left(A_{n}\right) x_{i}^{k_{i}+1}
$$

We express $p$ as

$$
p=\sum_{v \in V_{\ell}} c_{v} x^{v}
$$

for some $\ell$ such that $c_{v} \in \mathbb{Z}_{+}$. Then if $c_{v} \neq 0$, then $x^{v} \in J$. Hence, there must be a $v \in V_{\ell}$ such that $c_{v} \neq 0$ and $v_{i} \leq k_{i}$ for $i=1,2, \ldots, n-1$. Then $q=x_{0}^{v_{0}} x_{1}^{k_{1}} \cdots x_{n-1}^{k_{n-1}} \in$ $J \cap I(0, k)$ but $q \notin \sum_{i=1}^{n-1} K_{0}(I(0, k)) x_{i}$, which is a contradiction. This concludes the proof.

We note that $\pi_{k}$ is a $\sigma$-covariant representation. To show this define a unitary $U$ on $\ell^{2}\left(\Omega_{k}\right)$ by

$$
U \xi_{S}=\xi_{S+1}, \quad s \in \Omega_{k},
$$

where $S+1=\left(S_{0}+1, S_{1}+1, \ldots, S_{n-1}+1\right)$ for $S=\left(S_{0}, \ldots, S_{n-1}\right)$ and $S_{i}+1=$ $\left\{m+1 \mid m \in S_{i}\right\}$. Since

$$
U \pi_{k}(e(S, T)) U^{*} \xi_{V}=\delta_{T, V-1} U \xi_{S}=\delta_{T+1, V} \xi_{S+1}=\pi_{k}(e(S+1, T+1)) \xi_{V}
$$

for $S, T, V \in \Omega_{k}$, it follows that $\operatorname{Ad} U \pi_{k}(x)=\pi_{k} \sigma(x), x \in A_{n}$. If $|k|=1, \Omega_{k}$ is identified with $\mathbb{Z}$ and $U$ with the shift unitary $S$ on $\ell^{2}(\mathbb{Z})$. If $|k|>1$, then $U$ is unitarily equivalent to the shift unitary with infinite multiplicity.

From now on we denote by $\pi_{\left(\infty, k_{1}, \ldots, k_{n-1}\right)}$ the above representation $\pi_{k}$ with $k=$ $\left(k_{1}, \ldots, k_{n-1}\right)$. By assigning the role played by the index 0 to another index, we define an irreducible representation $\pi_{k}$ for $k \in\left(\mathbb{Z}_{+} \cup\{\infty\}\right)^{n}$ such that $k_{i}=\infty$ for a unique $i \in\{0,1, \ldots, n-1\}$. For such an index $k$, we define $\tilde{k} \in \mathbb{Z}_{+}^{n}$ by $\tilde{k}_{i}=k_{i}$ if $k_{i}<\infty$ and by $\tilde{k}_{i}=0$ if $k_{i}=\infty$. We then have that $\pi_{k}(I(\tilde{k}))=\mathcal{K}$ and

$$
\operatorname{Ker} \pi_{k}=\sum_{i=0}^{n-1} I\left(\left(k_{i}+1\right) e_{i}\right),
$$

where $I\left((\infty+1) e_{i}\right)=I\left(\infty e_{i}\right)=\{0\}$. Thus, we obtain the following lemma.
Lemma 3.2. Let $k \in\left(\mathbb{Z}_{+} \cup\{\infty\}\right)^{n}$ be such that $k_{i}=\infty$ for a unique $i$ and let $\tilde{k}$ be as above, and define a $\sigma$-covariant irreducible representation $\pi_{k}$ of $A_{n}$ with an implementing unitary $U$ as above. If $0<k_{j}<\infty$ for some $j, \pi_{k}(I(\tilde{k}))$ is isomorphic to the $C^{*}$-algebra of compact operators and $U$ is unitarily equivalent to $S \otimes 1$ on the Hilbert space $\ell^{2}(\mathbb{Z}) \otimes \mathcal{H}$, where $S$ is the shift unitary on $\ell^{2}(\mathbb{Z})$ and $\mathcal{H}$ is a separable Hilbert space, and $\mathcal{H}$ is one dimensional if $|\tilde{k}|=1$ and infinite dimensional if $|\tilde{k}|>1$.

In the following lemma, for a finite subset $X$ of $\mathbb{Z} \times \mathbb{Z}$, we denote by $P_{X}$ the projection onto the subspace generated by $\xi_{m}, m \in X$ in $\ell^{2}(\mathbb{Z} \times \mathbb{Z})=\ell^{2}(\mathbb{Z}) \otimes \ell^{2}(\mathbb{Z})$.

Lemma 3.3. Let $S$ be the shift unitary on $\ell^{2}(\mathbb{Z}) ; S \xi_{m}=\xi_{m+1}, m \in \mathbb{Z}$. For any finite subset $X$ of $\mathbb{Z} \times \mathbb{Z}$ and $\epsilon>0$, there is a finite subset $Y$ of $\mathbb{Z} \times \mathbb{Z}$ such that if $u \in \mathcal{K}\left(\ell^{2}(\mathbb{Z} \times \mathbb{Z})\right)+1$ with $u P_{Y}=P_{Y}$, then there is a unitary $v \in \mathcal{K}\left(\ell^{2}(\mathbb{Z} \times \mathbb{Z})\right)+1$ such that $v P_{X}=P_{X}$ and

$$
\left\|u-v(S \otimes 1) v^{*}\left(S^{*} \otimes 1\right)\right\|<\epsilon
$$

Moreover, if $X$ is empty, then $Y$ can also be set to be empty.
Proof. This is Lemma 2.4 of [19] when $X=\emptyset=Y$. Strengthening it to the above assertion is easy, but we present a proof below.

We may replace $\mathbb{Z} \otimes \mathbb{Z}$ by $\mathbb{Z} \otimes \Lambda$ with $\Lambda$ a finite non-empty set, and $X$ by $X^{\prime} \times \Lambda$ with $X^{\prime}$ a finite subset of $\mathbb{Z}$. We choose $Y$ as $Y^{\prime} \otimes \Lambda$ with $Y^{\prime}$ a finite subset of $\mathbb{Z}$.

We choose $Y^{\prime}$ so large that if $u$ is a unitary in $\mathcal{K}\left(\ell^{2}(\mathbb{Z} \times \Lambda)\right)+1$ with $u P_{Y}=P_{Y}$, then there is a projection $f \in \mathcal{K}\left(\ell^{2}(\mathbb{Z})\right)$ such that $\left\|S f S^{*}-f\right\|<\epsilon, f P_{X^{\prime}}=0$, and

$$
\|u-((f \otimes 1) u(f \otimes 1)+(1-f) \otimes 1)\|<\epsilon / 2
$$

There is a sequence $\left(g_{n}\right)$ of projections in $\mathcal{K}\left(\ell^{2}(\mathbb{Z})\right)$ such that $g_{n}\left(f+P_{X^{\prime}}\right)=0$, $\left\|S g_{n} S^{*}-g_{n}\right\| \rightarrow 0,\left\|u\left(g_{n} \otimes 1\right)-g_{n} \otimes 1\right\| \rightarrow 0, \operatorname{and} \operatorname{rank}\left(g_{n}\right) \rightarrow \infty$. Letting $f_{n}=f+g_{n}$, we have that $\left\|f_{n} S f_{n} S^{*} f_{n}-f_{n}\right\|<\epsilon$ for all sufficiently large $n$ and that

$$
\left\|\left(f_{n} \otimes 1\right) u(S \otimes 1)\left(f_{n} \otimes 1\right)\left(S^{*} \otimes 1\right) u^{*}\left(f_{n} \otimes 1\right)-f_{n} \otimes 1\right\|<2 \epsilon
$$

for all sufficiently large $n$. It follows that $f_{n} S f_{n} \otimes 1 \approx f S f^{*} \otimes 1+g_{n} S g_{n} \otimes 1$ and $\left(f_{n} \otimes 1\right) u(S \otimes 1)\left(f_{n} \otimes 1\right) \approx(f \otimes 1) u(S \otimes 1)(f \otimes 1)+g_{n} S g_{n} \otimes 1$ are approximated by unitaries on $f_{n} \ell^{2}(\mathbb{Z} \times \Lambda)$, up to the order $\epsilon$, whose spectra are almost uniformly distributed over $\mathbb{T}$ (due to the contribution from $g_{n} S g_{n} \otimes 1$ ). Hence, we find a unitary $v \in \mathcal{K}\left(\ell^{2}(\mathbb{Z} \times \Lambda)\right)+1$ such that $v=v\left(f_{n} \otimes 1\right)+\left(1-f_{n}\right) \otimes 1$ for a sufficiently large $n$ and $\left\|\left(f_{n} \otimes 1\right) u(S \otimes 1)\left(f_{n} \otimes 1\right)-v\left(f_{n} \otimes 1\right)(S \otimes 1)\left(f_{n} \otimes 1\right) v^{*}\right\|$ is at most of order $\epsilon$. This implies that $\left\|u(S \otimes 1)-v(S \otimes 1) v^{*}\right\|$ is at most of order $\epsilon$, concluding the proof.

## 4. The one-cocycle and the approximate Rohlin properties

For a subset $F$ of $\mathbb{Z}$ we set

$$
A_{n}(F)=\left(\bigotimes_{F} M_{n}\right)^{\gamma}
$$

as a $C^{*}$-subalgebra of $A_{n}$. If $m \in \mathbb{Z}_{+}$and $k \in \mathbb{N}^{m}$, let $Q$ be the quotient map of $A_{n+m}$ onto

$$
B=A_{n+m} / \sum_{i=0}^{m-1} I\left(k_{i} e_{n+i}\right)
$$

For a subset $F$ of $\mathbb{Z}$ we set $B(F)=Q\left(A_{n+m}(F)\right)$, i.e. the local structure of $B$ is defined by $\{B(F) \mid F \subset \mathbb{Z}\}$.

Definition 4.1. Let $n=2,3, \ldots$. We say that the $n$-shift has the one-cocycle property if the following conditions are satisfied.
(1) For any $\epsilon>0$ and any unitary $u \in A_{n}$ such that $\varphi_{i}(u)=1$ for $i=0, \ldots, n-1$, there is a unitary $v \in A_{n}$ such that $\left\|u-v \sigma(v)^{*}\right\|<\epsilon$.
(2) If ( $u_{k}$ ) is a sequence in the unitary group $\mathcal{U}\left(A_{n}\right)$ and $\left(F_{k}\right)$ is an increasing sequence of finite subsets of $\mathbb{Z}$ such that $\varphi_{i}\left(u_{k}\right)=1$ for $i=0, \ldots, n-1, \bigcup_{k} F_{k}=\mathbb{Z}$, and $u_{k} \in A_{n}\left(\mathbb{Z} \backslash F_{k}\right)$, then there is a sequence $\left(v_{k}\right)$ in $\mathcal{U}\left(A_{n}\right)$ and an increasing sequence $\left(G_{k}\right)$ of finite subsets of $\mathbb{Z}$ such that $\left\|u_{k}-v_{k} \sigma\left(v_{k}\right)\right\| \rightarrow 0, \bigcup_{k} G_{k}=\mathbb{Z}$, and $v_{k} \in A_{n}\left(\mathbb{Z} \backslash G_{k}\right)$.
(3) If $\left(u_{k}\right)$ is a sequence in $\mathcal{U}\left(A_{n}\right)$ such that $\varphi_{i}\left(u_{k}\right)=1$ for $i=0, \ldots, n-1$ and $\left\|\left[u_{n}, x\right]\right\| \rightarrow 0$ for all $x \in A_{n}$, then there is a sequence $\left(v_{k}\right)$ in $\mathcal{U}\left(A_{n}\right)$ such that $\left\|u_{k}-v_{k} \sigma\left(v_{k}\right)\right\| \rightarrow 0$ and $\left\|\left[v_{k}, x\right]\right\| \rightarrow 0$ for all $x \in A_{n}$.
We say that the $n$-shift has the stable one-cocycle property if the above conditions are satisfied for the shift $\sigma$ on the quotient $B=A_{n+m} / \sum_{i=0}^{m-1} I\left(k_{i} e_{n+i}\right)$ for any $m \in \mathbb{Z}_{+}$and $k \in \mathbb{N}^{m}$.

By the definition if the $n$-shift has the stable one-cocycle property then it has the onecocycle property. If the $n$-shift has the one-cocycle property for all $n \geq 2$, then the $n$-shift has the stable one-cocycle property. Since each of the properties in Theorem 1.1 is equivalent to the corresponding one of the above definition, the main Theorem 1.1 states that the $n$-shift has the stable one-cocycle property for all $n$.

Let $\operatorname{Ex}\left(\Lambda_{n}\right)$ be the set of extreme points of $\Lambda_{n}$, i.e. $\operatorname{Ex}\left(\Lambda_{n}\right)$ consists of $n$ vertices. We call an $N$-cycle an orthogonal family of projections indexed by $\mathbb{Z} / N \mathbb{Z}$.
Definition 4.2. Let $n=2,3, \ldots$. We say that the $n$-shift has the approximate Rohlin property if for any $N \in \mathbb{N}$ and $\epsilon>0$, there exist an $N_{1}, S \in \mathbb{N}$ with $N_{1} \geq N$ and a constant $C>0$ such that there is a sequence $\left(e_{i}^{(k)}\right)$ of $N_{1}$-cycles in $A_{n}$ with $E_{k}=\sum_{i} e_{i}^{(k)}$ satisfying the following:

$$
\begin{aligned}
& \max _{i}\left\|\sigma\left(e_{i}^{(k)}\right)-e_{i+1}^{(k)}\right\| \rightarrow 0, \quad k \rightarrow \infty \\
& {\left[E_{k}\right](x) \geq C\left(x_{0} x_{1} \cdots x_{n-1}\right)^{S}, \quad x \in \Lambda_{n}} \\
& {\left[E_{k}\right](x) \geq c_{k} \sum_{i \neq j}\left(x_{i} x_{j}\right)^{S}, \quad x \in \Lambda_{n}}
\end{aligned}
$$

for some $c_{k}>0$, and $\left[E_{k}\right](x)$ converges to 1 uniformly in $x$ on every compact subset of $\Lambda_{n} \backslash \operatorname{Ex}\left(\Lambda_{n}\right)$.

Note that when we define a Rohlin property we usually impose a centrality condition such as $\left\|\left[e_{i}^{(k)}, x\right]\right\| \rightarrow 0$ as $k \rightarrow \infty$ for all $x$. Since we are dealing with the shift here, such a condition follows automatically.

To see that the lower bound estimates on $\left[E_{k}\right]$ are not redundant, we can construct an automorphism $\alpha$ of $A_{n}$ which has an approximate Rohlin property without the lower bound estimates $\left[E_{k}\right]$ above, but with a proper centrality condition. For example, we choose a sequence $\left(F_{k}\right)$ of finite subsets of $\mathbb{Z}$ which are mutually disjoint such that $\left|F_{k}\right| / k \rightarrow \infty$ and choose a sequence $\left(U_{k}\right)$ of unitaries in $A_{n}$ such that $U_{k} \in A_{n}\left(F_{k}\right)$ and $U_{k} E_{v}=E_{v}$ for $v \in V_{\left|F_{k}\right|}$ with $\min _{i} v_{i}<k$ and the spectrum of $U_{k} E_{v}$ is equally distributed otherwise, where $E_{v}$ is the minimal central projection of $A_{n}\left(F_{k}\right)$ as in the proof of Proposition 2.1. And we define an automorphism $\alpha$ of $A_{n}$ by the limit of $\operatorname{Ad}\left(U_{1} U_{2} \cdots U_{k}\right)$. Then we can see that $\alpha$ has the property that for any $N \in \mathbb{N}$ there is a sequence $\left(e_{i}^{(k)}\right)$ of $N$-cycles in $A_{n}$ such that $\max _{i}\left\|\alpha\left(e_{i}^{(k)}\right)-e_{i+1}^{(k)}\right\| \rightarrow 0, N\left[e_{0}^{(k)}\right](x) \rightarrow 1$ uniformly in $x$ on every compact
subset of $\Lambda_{n} \backslash \operatorname{Ex}\left(\Lambda_{n}\right)$, and $\left\|\left[e_{i}^{(k)}, x\right]\right\| \rightarrow 0$ for $x \in A_{n}$ (see, e.g., [2]). However, then the ideal $J_{k}$ of $A_{n}$ generated by $e_{i}^{(k)}$ must satisfy $\bigcap_{k} J_{k}=\{0\}$.

Our purpose is to prove, by induction, that the properties defined above follow universally. The following lemmas give the basic step for this induction, which are more or less presented in [19].

Lemma 4.3. The 2 -shift has the approximate Rohlin property.
Proof. By using the CAR algebra formulation we show in [19] that for any $N \in \mathbb{N}$ there is a sequence of $2^{N}$-cycles in $A_{2}$ satisfying the required properties with $M=N$ and $C=2^{N}$.

Lemma 4.4. The 1 -shift has the stable one-cocycle property.
Proof. Let $m \in \mathbb{N}$ and $k \in \mathbb{Z}_{+}^{1+m}$ such that $k_{0}=0$ and $k_{j}>0$ for $j>0$ and define

$$
B=A_{1+m} / \sum_{i=1}^{m} I\left(k_{i} e_{i}\right)
$$

We have to show that the shift $\sigma$ defined on $B$ has the one-cocycle property.
Let $\varphi_{0}$ be the unique character of $B$. What we have to prove is that for any $u \in \mathcal{U}(B)$ with $\varphi_{0}(u)=1$, there is a sequence $\left(v_{n}\right)$ in $\mathcal{U}(B)$ such that $\left\|u-v_{n} \sigma\left(v_{n}\right)^{*}\right\| \rightarrow 0$; and two other versions.

By the following Lemma 4.5 and Lemmas 3.2 and 3.3 we find a finite decreasing sequence $\left(J_{i}\right)_{i=0}^{N}$ of ideals of $A_{1+m}$ such that $J_{0}=A_{1+m}, \quad J_{1}=\sum_{i=1}^{m} I\left(e_{i}\right)$, $J_{N} \subset \sum_{i=1}^{m} I\left(k_{i} e_{i}\right), J_{N-1} \not \subset \sum_{j=1}^{m} I\left(k_{j} e_{j}\right), J_{i-1} / J_{i} \cong \mathcal{K}$ for $i>1$, and the automorphism, denoted by $\sigma$, induced on $J_{i-1} / J_{i}$ by the shift $\sigma$ has the one-cocycle property. (In Lemma 4.5, since $J_{N}=\operatorname{Ker}\left(\rho_{N}\right) \cap J_{N-1}=\sum_{j=1}^{m} I\left(k_{j} e_{j}\right) \cap J_{N-1}$, it follows that $J_{N-1} \not \subset \sum_{j=1}^{m} I\left(k_{j} e_{j}\right)$.)

Let $u_{1} \in \mathcal{U}\left(A_{1+m}\right)$ be such that $Q\left(u_{1}\right)=u$, where $Q$ is the quotient map of $A_{1+m}$ onto $B$. Since $\varphi_{0}\left(Q\left(u_{1}\right)\right)=1$, we have that $u_{1} \in J_{1}+1$.

From the short exact sequence

$$
0 \rightarrow J_{2} \rightarrow J_{1} \rightarrow J_{1} / J_{2} \rightarrow 0
$$

and the one-cocycle property for the shift on $J_{1} / J_{2}$ applied to $u_{1}+J_{2}$, proved in Lemma 3.3, we find unitaries $u_{2} \in J_{2}+1$ and $v_{1} \in J_{1}+1$ such that

$$
\left\|v_{1}^{*} u_{1} \sigma\left(v_{1}\right)-u_{2}\right\|<N^{-1} \epsilon
$$

Repeating this process, we have unitaries $u_{i} \in J_{i}+1$ for $i \leq N$ and unitaries $v_{i} \in J_{i}+1$ for $i<N$ such that

$$
\left\|v_{i}^{*} u_{i} \sigma\left(v_{i}\right)-u_{i+1}\right\|<N^{-1} \epsilon
$$

for $i<N$. Then it follows that

$$
\left\|v_{1} v_{2} \cdots v_{N-1} u_{N} \sigma\left(v_{1} v_{2} \cdots v_{N-1}\right)-u_{1}\right\|<\epsilon
$$

Thus, for $v=Q\left(v_{1} v_{2} \cdots v_{N-1}\right) \in B$ we get that $\left\|v \sigma(v)^{*}-u\right\|<\epsilon$. This concludes the proof of the first property of Definition 4.1.

We turn to the second property of Definition 4.1. Let $G$ be a finite subset of $\mathbb{Z}$. To get $v_{N-1}$ from $A_{1+m}(\mathbb{Z} \backslash G) \cap J_{N-1}+1$ in the above arguments based on Lemma 3.3, $u_{N-1}$ must be from $A_{1+m}(\mathbb{Z} \backslash F)$ for some finite subset $F_{1} \subset \mathbb{Z}$ with $F_{1} \supset G$. Here we used the fact that requiring $v_{N-1} \in A_{1+m}(\mathbb{Z} \backslash G) \cap J_{N-1}+1$ modulo $J_{N}$ is as strong as requiring that $v_{N-1} p \in \mathbb{C} p$ for some $p \in J_{N-1}$ which is a projection modulo $J_{N}$ since $J_{N-1} / J_{N} \cong \mathcal{K}$. For that to be true, we must have that $v_{N-2}, u_{N-2} \in A_{1+m}\left(\mathbb{Z} \backslash F_{1}\right)$, which in turn requires that $u_{N-2} \in A_{1+m}\left(\mathbb{Z} \backslash F_{2}\right)$ for some finite subset $F_{2} \supset F_{1}$. In this way we can conclude that there is a finite subset $F_{N-1}$ of $\mathbb{Z}$ such that if $u_{1}$ is a unitary from $A_{1+m}\left(\mathbb{Z} \backslash F_{N-1}\right)$, then all $v_{1}, \ldots, v_{N-1}$ will be chosen from $A_{1+m}(\mathbb{Z} \backslash G)$ in the above arguments; then $v=Q\left(v_{1} \cdots v_{N-1}\right) \in B(\mathbb{Z} \backslash G)$ is the desired unitary. This concludes the proof in this case.

Finally, we come to the last property of Definition 4.1. Let $\mathcal{G}$ be a finite subset of $B$; by lifting each element of $\mathcal{G}$ to $A_{1+m}$ we regard $\mathcal{G}$ as a subset of $A_{1+m}$. We choose a finite subset $G$ of $\mathbb{Z}$ such that if $v$ is a unitary in $A_{1+m}(\mathbb{Z} \backslash G)$ then $\|[v, x]\|<1 / 2 N$, $x \in \mathcal{G}$. To get $v_{N-1} \in J_{N-1}+1$ for $u_{N-1}$ such that $v_{N-1} \in A_{1+m}(\mathbb{Z} \backslash G)$, we must require that $\left\|\left[u_{N-1}, x\right]+J_{N}\right\|<1, x \in \mathcal{F}_{1}$ for some finite subset of $\mathcal{F}_{1}$ of $A_{1+m}$. For that we have to require that $v_{N-2} \in A_{1+m}\left(\mathbb{Z} \backslash G_{1}\right)$ and $\left\|\left[u_{N-2}, x\right]+J_{N}\right\|<1 / 3 N, x \in \mathcal{F}_{1}$ for some finite subset $G_{1}$ of $\mathbb{Z}$ with $G_{1} \supset G$. Again for this to be true, we have to require that $\left\|\left[u_{N-2}, x\right]+J_{N}\right\|<1 / N, x \in \mathcal{F}_{2}$ for some finite subset $\mathcal{F}_{2}$ of $A_{1+m}$. In this way we obtain a finite subset $\mathcal{F}_{N-1}$ of $A_{1+m}$ such that if $u_{1} \in J_{1}+1$ satisfies that $\left\|\left[u_{1}, x\right]+J_{N}\right\|<1 / N$ for $x \in \mathcal{F}_{N-1}$, then all $v_{1}, \ldots, v_{N-1}$ can be taken from $A_{1+m}(\mathbb{Z} \backslash G)$. Then $v=Q\left(v_{1} \cdots v_{N-1}\right)$ would satisfy the required condition. (Thus, since $B=A_{1+m} / \sum_{j=1}^{m} I\left(k_{j} e_{j}\right)$ is of type I, properties (2) and (3) are equivalent for this case.)

LEMMA 4.5. There exists a finite sequence $\left(\rho_{i}\right)_{i=1}^{N}$ of irreducible representations of $A_{1+m}$ of the form $\pi_{(\infty, \ell)}, \ell \in\left(\mathbb{Z}_{+}\right)^{m}$ and a decreasing sequence $\left(J_{i}\right)_{i=0}^{N}$ of ideals of $A_{1+m}$ such that $J_{0}=A_{1+m}, \rho\left(J_{i-1}\right)=\mathcal{K}$ or $\mathbb{C}, \operatorname{Ker}\left(\rho_{i} \mid J_{i-1}\right)=J_{i}$, and

$$
\begin{aligned}
\operatorname{Ker}\left(\rho_{i}\right) & \supset \sum_{j=1}^{m} I\left(k_{j} e_{j}\right), \\
J_{N} & \subset \sum_{j=1}^{m} I\left(k_{j} e_{j}\right) .
\end{aligned}
$$

Proof. Let $\rho_{1}=\pi_{(\infty, 0, \ldots, 0)}$ and $J_{1}=\operatorname{Ker} \rho_{1}$. Then it follows that $J_{1}=\sum_{j=1}^{m} I\left(e_{j}\right) \supset$ $\sum_{j=1}^{m} I\left(k_{j} e_{j}\right)$. We prove the assertion by induction as follows.

Suppose that we are given an ideal $J$ of $A_{1+m}$ such that

$$
J=\sum_{w \in \mathcal{C}} I(w)
$$

where $\mathcal{C}$ is a subset of $\mathbb{Z}_{+}^{1+m}$ such that $w \leq k$ and $|w|=\sum_{j=0}^{m} w_{j} \geq 1$ for $w \in \mathcal{C}$. We may suppose that no pairs in $\mathcal{C}$ are comparable. (If $w \geq w^{\prime}$ in the sense that $w_{j} \geq w_{j}^{\prime}$, then $I(w) \subset I\left(w^{\prime}\right)$; so we may remove $w$ from $\mathcal{C}$.)

Let

$$
r=\min \left\{|w| \mid w \in \mathcal{C}, w_{j}<k_{j} \text { for all } j=1, \ldots, m\right\} .
$$

For each $w \in\left\{w \in \mathcal{C}||w|=r\}\right.$, we define $\hat{w} \in\left(\mathbb{Z}_{+} \cup\{\infty\}\right)^{1+m}$ by $\hat{w}_{0}=\infty$ and $\hat{w}_{j}=w_{j}$ for $j>0$ and consider $\pi_{\hat{w}}$. For $z \in \mathcal{C} \backslash\{w\}$, we have that $I(z) \subset \operatorname{Ker} \pi_{\hat{w}}$. (Either $z_{j} \geq k_{j}>w_{j}$ for some $j=1,2, \ldots, m$ or $z_{j}>w_{j}$ for some $j>0$ as $|z| \geq|w|$.) Then $\pi_{\hat{w}}(J)=\mathcal{K}$ and

$$
\operatorname{Ker}\left(\pi_{\hat{w}} \mid J\right)=\sum_{z \in \mathcal{C} \backslash\{w\}} I(z)+\sum_{j=1}^{m} I\left(w+e_{j}\right)
$$

Note also that $\operatorname{Ker} \pi_{\hat{w}} \supset \sum_{j=1}^{m} I\left(\left(w_{j}+1\right) e_{j}\right) \supset \sum_{j=1}^{m} I\left(k_{j} e_{j}\right)$.
We now set $J$ to be $\operatorname{Ker}\left(\pi_{\hat{w}} \mid J\right)$ and $\mathcal{C}$ to be an appropriate subset of $\mathcal{C} \backslash\{w\} \cup\left\{w+e_{j} \mid\right.$ $j=1, \ldots, m\}$. Repeating this process for a finite number of times we can increase $r=\min \left\{|w| \mid w \in \mathcal{C}, w_{j}<k_{j}\right.$ for all $\left.j>0\right\}$. Eventually we reach the situation where for any $w \in \mathcal{C}$ there is $j>0$ such that $w_{j}=k_{j}$. Hence, we get $J=\sum_{w \in \mathcal{C}} I(w) \subset$ $\sum_{j=1}^{m} I\left(k_{j} e_{j}\right)$. This completes the proof.

The following result combined with Lemma 4.3 shows that the $n$-shift has the approximate Rohlin property for all $n \geq 2$.

Lemma 4.6. Let $n=2,3, \ldots$. If the $n$-shift has the approximate Rohlin property, then the $n+1$-shift has also the approximate Rohlin property.
Proof. Let $K, L, M \in \mathbb{N}$ be such that $M \gg L$ and let $Z_{M, K}=\{m M \mid m=0,1, \ldots$, $K-1\}$ and $Z_{M, K, L}=\bigcup\left\{Z_{M, K}+m \mid m=0,1, \ldots, L-1\right\}$. We define an embedding $\iota$ of $M_{n}$ into $\bigotimes_{\{0\}} M_{n+1} \subset \bigotimes_{\mathbb{Z}} M_{n+1}$ by $\iota\left(e_{i j}\right)=e_{i j}^{(0)}$, where $\left(e_{i j}^{(0)}\right)$ are the canonical matrix units for $M_{n+1}$ at $0 \in \mathbb{Z}$ as a $C^{*}$-subalgebra of $\bigotimes_{\mathbb{Z}} M_{n+1}$; in particular, $1-\iota_{0}(1)=e_{n n}^{(0)}$. We then define a homomorphism $\psi_{0}$ of $M_{n}$ into $\bigotimes_{Z_{M, K}} M_{n+1} \subset \bigotimes_{\mathbb{Z}} M_{n+1}$ by

$$
\begin{aligned}
\psi_{0}(x)= & \iota_{0}(x)+e_{n n}^{0} \sigma^{M}\left(\iota_{0}(x)\right)+e_{n n}^{(0)} e_{n n}^{(L)} \sigma^{2 M}\left(\iota_{0}(x)\right)+\cdots \\
& +e_{n n}^{(0)} \cdots e_{n n}^{((K-2) L)} \sigma^{(K-1) M}\left(\iota_{0}(x)\right) .
\end{aligned}
$$

We set $\psi_{k}=\sigma^{k} \psi_{0}$ for $k=1,2, \ldots, M-1$; note that the ranges of $\psi_{k}, k=0,1, \ldots, M-1$ mutually commute and that $L \ll M$. We then define a homomorphism $\psi$ of $\bigotimes_{0}^{L-1} M_{n}$ into $\otimes_{Z_{M, K, L}} M_{n+1} \subset \otimes_{\mathbb{Z}} M_{n+1}$ by $\psi=\psi_{0} \otimes \psi_{1} \otimes \cdots \otimes \psi_{L-1}$. Note that $\psi_{L}(1) \psi \sigma \stackrel{Z_{M, K, L}}{=} \psi_{0}(1) \sigma \psi$ on $\bigotimes_{0}^{L-2} M_{n}$ and that the range of $\psi \mid\left(\bigotimes_{0}^{L-1} M_{n}\right)^{\gamma}$ is contained in $A_{n+1}=\left(\otimes_{\mathbb{Z}} M_{n+1}\right)^{\gamma}$, where $\gamma$ is the gauge action of $\mathbb{T}^{n-1}$ (or $\mathbb{T}^{n}$ for the latter case). We thus regard $\psi$ as a homomorphism of $B_{L}=\left(\otimes_{0}^{L-1} M_{n}\right)^{\gamma}$ into $A_{n+1}$.

We recall that there is a tracial state $\tau_{\lambda}$ on $A_{n+1}$ to each

$$
\lambda \in \Lambda_{n+1}=\left\{\lambda \in[0,1]^{n} \mid \sum_{i=1}^{n} \lambda_{i} \leq 1\right\} .
$$

Note that $\tau_{\lambda} \psi_{0}(1)=1-\lambda_{n}^{K}$ and $\tau_{\lambda} \psi(1)=\left(1-\lambda_{n}^{K}\right)^{L}$. For a projection $e \in B_{L}$, we have that $\tau_{\lambda} \psi(e)=0$ for $\lambda_{n}=1$ and

$$
\tau_{\lambda} \psi(e)=\left(1-\lambda_{n}^{K}\right)^{L} \tau_{\mu}(e)
$$

for $\lambda_{n}<1$, where

$$
\mu=\left(\frac{\lambda_{1}}{1-\lambda_{n}}, \ldots, \frac{\lambda_{n-1}}{1-\lambda_{n}}\right) \in \Lambda_{n}
$$

Hence, if $[e]=g \in \mathbb{Z}\left[y_{1}, \ldots, y_{n-1}\right]$ in $K_{0}\left(A_{n}\right)$ under the embedding $B_{L} \subset A_{n}$, then $f=[\psi(e)]$ in $K_{0}\left(A_{n+1}\right)$ is given by

$$
f\left(x_{1}, \ldots, x_{n}\right)=\left(1-x_{n}^{K}\right)^{L} g\left(x_{1} /\left(1-x_{n}\right), \ldots, x_{n-1} /\left(1-x_{n}\right)\right),
$$

which must be a polynomial in $x_{1}, \ldots, x_{n}$, since $g$ is at most of order $L$. If $g$ is strictly positive on $\Lambda_{n} \backslash \operatorname{Ex}\left(\Lambda_{n}\right)$, then $f$ is strictly positive on $\Lambda_{n+1} \backslash \bigcup_{i=0}^{n-1}[i, n]$, where $[i, n]$ is the edges between $x_{i}=1$ and $x_{n}=1$ for $i=0,1, \ldots, n-1$. We can make $f$ arbitrarily close to 1 uniformly on a compact subset of $\Lambda_{n+1} \backslash \bigcup_{i=0}^{n-1}[i, n]$ by making $g$ close to 1 on a compact subset of $\Lambda_{n} \backslash \operatorname{Ex}\left(\Lambda_{n}\right)$. If $g\left(y_{1}, \ldots, y_{n-1}\right) \geq C\left(y_{0} \cdots y_{n-1}\right)^{S}$ for some $S \in \mathbb{N}$ and $C>0$, where $y_{0}=1-y_{1}-\cdots-y_{n-1}$, then $f(x) \geq C\left(1-x_{n}^{K}\right)^{L}\left(1-x_{n}\right)^{-n S}\left(x_{0} \cdots x_{n-1}\right)^{S}$ for $x \in \Lambda_{n+1}$ because

$$
1-\frac{x_{1}}{1-x_{n}}-\frac{x_{2}}{1-x_{n}}-\cdots-\frac{x_{n-1}}{1-x_{n}}=\frac{x_{0}}{1-x_{n}}
$$

Note that if $K$ is sufficiently large, then $\left(1-x_{n}^{K}\right)^{L} \approx 1$ uniformly for $x_{n}$ in [0,1- $\delta$ ] for $\delta>0$. Thus, for such a choice of $K$, we have that

$$
f(x) \geq C\left(x_{0} x_{1} \cdots x_{n-1}\right)^{S}
$$

for $x \in \Lambda_{n+1}$ with $x_{n} \in[0,1-\delta]$.
Since $\psi_{L}(1) \psi \sigma=\psi_{0}(1) \sigma \psi$ on $\left(\otimes_{0}^{L-2} M_{n}\right)^{\gamma} \otimes 1 \subset B_{L}, \psi$ does not really intertwine the shifts $\sigma$ even approximately. We have to modify $\psi$ as follows.

By using the lemma below and taking a sufficiently large $M$, we define an embedding $\phi: B_{L} \rightarrow A_{n+1}$ by $\phi(x)=\psi(x) p$, where $p$ is a projection which resides outside $Z_{M, K, L}$ and satisfies that $\|\sigma(\phi(1))-\phi(1)\|<\epsilon$ for a prescribed $\epsilon>0$ and $\tau_{\lambda}(\phi(1)) \geq\left(1-\lambda_{n}^{K}\right)^{L+L_{\epsilon}}$ with some $L_{\epsilon} \in \mathbb{N}$. Since $\|\sigma(\psi(1) p)-\psi(1) p\|<\epsilon$ and $\sigma(\psi(1) p)=\psi_{L}(1) \sigma(\psi(1) p)$, it follows that $\left\|\psi_{L}(1) \psi(1) p-\psi(1) p\right\|<2 \epsilon$. If $e \in$ $B_{L} \subset A_{n}$ is a projection such that $\sigma(e) \in B_{L}$, then it follows that $\|\sigma \phi(e)-\phi \sigma(e)\|<$ $3 \epsilon$ because $\sigma \phi(e)=\sigma(\psi(e)) \sigma(\psi(1) p) \approx \sigma \psi(e) \psi(1) p=\sigma \psi(e) \psi_{0}(1) \psi(1) p=$ $\psi \sigma(e) \psi_{L}(1) \psi(1) p \approx \psi \sigma(e) p=\phi \sigma(e)$. Since this modification of $\psi$ introduces only the factor $\left(1-\lambda_{n}^{K}\right)^{L_{\epsilon}}$ in the estimate of $\tau_{\lambda}(\phi(e))$ for a projection $e \in B_{L}$, we have a similar estimate for $f=[\phi(e)]$ as in the previous paragraph.

Let $N \in \mathbb{N}$. Then there are $N_{1}, C, S \in \mathbb{N}$ with $N_{1} \geq N$ satisfying the following: we have an $N_{1}$-cycle $\left(e_{i}\right)$ in $A_{n}$ such that $\left\|\sigma\left(e_{i}\right)-e_{i+1}\right\| \approx 0, N_{1}\left[e_{i}\right]\left(x_{1}, \ldots, x_{n-1}\right)$ is close to 1 uniformly on a compact subset of $\Lambda_{n} \backslash \operatorname{Ex}\left(\Lambda_{n}\right)$, and

$$
\left[e_{i}\right]\left(x_{1}, \ldots, x_{n-1}\right) \geq C\left(x_{0} x_{1} \cdots x_{n-1}\right)^{S}
$$

on $\Lambda_{n}$, and

$$
\left[e_{i}\right]\left(x_{1}, \ldots, x_{n-1}\right) \geq c \sum_{i \neq j}\left(x_{i} x_{j}\right)^{S}
$$

on $\Lambda_{n}$ for some $c>0$ (which may depend on the choice of $e_{i}$ ).
We may assume that $e_{i} \in \bigotimes_{m=0}^{L-1} M_{n}$ for some $L$ and we embed $e_{i}$ into $A_{n+1}$ by $\phi$ which depends on the choice of $K, M, p$. Thus, for any $\epsilon>0$ we get an $N_{1}$-cycle ( $E_{i}$ ) in $A_{n+1}$ such that $\left\|\sigma\left(E_{i}\right)-E_{i+1}\right\|<\epsilon$,

$$
N_{1}\left[E_{i}\right]\left(x_{1}, \ldots, x_{n}\right)>1-\epsilon
$$

for $x \in \Lambda_{n+1}$ such that $\min _{i \neq n} \operatorname{dist}(x,[i, n])>\epsilon$, where the distance dist on $\Lambda_{n+1}$ is defined by $\operatorname{dist}(\lambda, \nu)=\max \left\{\left|\lambda_{i}-\mu_{i}\right| ; i=0,1, \ldots, n\right\}$, and

$$
\left[E_{i}\right]\left(x_{1}, \ldots, x_{n}\right) \geq C\left(x_{0} x_{1} \cdots x_{n}\right)^{S}
$$

on the subset of $\Lambda_{n+1}$ consisting of $x$ with $x_{n}<1-\epsilon$, and

$$
\left[E_{i}\right]\left(x_{1}, \ldots, x_{n}\right) \geq c \sum_{i \neq j \leq n-1}\left(x_{i} x_{j}\right)^{S}
$$

on the above subset for some $c>0$.
In this way we embed a Rohlin cycle in $A_{n}$ into $A_{n+1}$. Note that there are $n+1$ types of embedding by assigning the role played by the $n$th coordinate in the above to the other coordinates. We combine the Rohlin cycles so obtained to get the desired one. The proof is continued after the following lemmas.

Lemma 4.7. For any $\epsilon>0$ there exists an $L_{\epsilon} \in \mathbb{N}$ such that if $M>L+2 L_{\epsilon}$, there is a projection $p$ in $\left(\otimes_{Z^{\prime}} M_{n+1}\right)^{\gamma}$, where $Z^{\prime}=\bigcup\left\{Z_{M, K}+m \mid m=-1,-2, \ldots,-L_{\epsilon}\right.$; $\left.m=L, L+1, \ldots, L+L_{\epsilon}-1\right\}$ is disjoint from $Z_{M, K, L}=\bigcup\left\{Z_{M, K}+m \mid m=\right.$ $0,1, \ldots, L-1\}$ such that $f=\psi(1) p$ is a projection satisfying $\|\sigma(f)-f\|<\epsilon$ and $\tau_{\lambda}(f) \geq\left(1-\lambda_{n}^{K}\right)^{L+2 L_{\epsilon}}$ for $\lambda \in \Lambda_{n+1}$.

Proof. Let $e=\psi_{0}(1)$ as a projection in $\left(\otimes_{Z_{M, K}} M_{n+1}\right)^{\gamma}$. Then $\psi(1)$ looks like $e \otimes e \otimes \cdots \otimes e\left(\operatorname{in}\left(\otimes_{Z_{M, K, L}} M_{n+1}\right)^{\gamma}=\left(\otimes_{0}^{L-1} \otimes_{Z_{M, K}} M_{n+1}\right)^{\gamma}\right)$, where $e$ repeats $L$ times. If $Z^{\prime \prime}=Z^{\prime} \cup Z_{M, K, L}$, then $\bigotimes_{Z^{\prime \prime}} M_{n+1}$ is the tensor product of $\bigotimes_{Z_{M, K}} M_{n+1}$ indexed by $m$ between $-L_{\epsilon}$ and $L+L_{\epsilon}-1$ inclusive, where $\sigma$ shifts each factor to right and $\psi(1)$ resides at the tensor product between 0 and $L-1$. Hence, the problem reduces to finding an almost $\sigma$ invariant projection in $\left(\otimes_{-L_{\epsilon}}^{L+l_{\epsilon}-2} M_{n+1}\right)^{\gamma}$ dominates $e \otimes e \otimes \cdots \otimes e$ ( $L+2 L_{\epsilon}$ times) and is dominated by $\psi(1)=e \otimes e \otimes \cdots \otimes e$ ( $L$ times). This follows from the following lemma.

Lemma 4.8. Let $m \in \mathbb{N}$ and let $A=\left(\otimes_{\mathbb{Z}} M_{m}\right)^{\gamma}$, where $\gamma$ is a restriction of the (infinite tensor product type) action $\beta$ of $\mathcal{U}_{m}=\mathcal{U}\left(M_{m}\right)$ to a compact abelian subgroup of $\mathcal{U}_{m}$ and let $\sigma$ be the shift automorphism of $A$ to the right. Let $e$ be a projection in $\left(M_{m}\right)^{\gamma}$. For $K \in \mathbb{N}$, let $e_{K}=\bigotimes_{-K}^{K}$ e. For $L, N \in \mathbb{N}$ there is a projection $p \in\left(\bigotimes_{-N-2 L}^{N+2 L} M_{m}\right)^{\gamma} \subset A$ such that $e_{N+L} \leq p \leq e_{N}$ and $\|\sigma(p)-p\|$ is of the order of $L^{-1 / 2}$. Moreover, $p$ can be chosen to satisfy that if $\varphi$ is a state on $M_{m}$ and $\tau=\bigotimes_{\mathbb{Z}} \varphi$ is a tracial state on $A$, then

$$
\tau(p)=\varphi(e)^{2 N+2 L+1}(1+L(1-\varphi(e))) .
$$

Proof. Let $f=e_{N+L}$ and let $f_{k}=\sigma^{k}(f)\left(1-\sigma^{k-1}(f)\right)$ and $f_{-k}=\sigma^{-k}(f)(1-$ $\left.\sigma^{-k+1}(f)\right)$ for $k=1,2, \ldots, L$, which are all projections. If $1 \leq k<\ell$, then

$$
f_{k} f_{\ell} \leq \sigma^{k}(f) \sigma^{\ell}(f)\left(1-\sigma^{\ell-1}(f)\right)=0
$$

because $\sigma^{k}(f) \sigma^{\ell-1}(f) \geq \sigma^{k}(f) \sigma^{\ell}(f)$. If $k, \ell \geq 1$, then

$$
f_{-k} f_{\ell} \leq \sigma^{-k}(f)\left(1-\sigma^{-k+1}(f)\right) \sigma^{\ell}(f)=0
$$

because $\sigma^{-k}(f) \sigma^{\ell}(f) \leq \sigma^{k}(f) \sigma^{-k+1}(f)$. Thus, one can conclude that the projections $f_{k}, 0<|k| \leq L$ are mutually orthogonal and orthogonal to $f$. Note also that they are mutually equivalent and that $\sigma$ shifts the sequence of projections

$$
f_{-L}, f_{-L+1}, \ldots, f_{-1}
$$

to the right, the last one $f_{-1}$ to $f(1-\sigma(f))$, which is a subprojection of $f$, and a subprojection, $f\left(1-\sigma^{-1}(f)\right)$, of $f$ to $f_{1}$, and shifts the sequence

$$
f_{1}, f_{2}, \ldots, f_{L}
$$

to the right except for the last $f_{L}$. Let $v \in A$ be such that $v^{*} v=f_{-L}$ and $v v^{*}=f_{1}$ and define

$$
p=f+\sum_{k=1}^{L}\left(\frac{k}{L} f_{-L+k-1}+\frac{L-k}{L} f_{k}+\frac{\sqrt{k(L-k)}}{L} \sigma^{k-1}\left(v+v^{*}\right)\right) .
$$

Since all $f_{k}, f_{-k}, f$ are dominated by $e_{N}$ and $p \geq f=e_{N+L}$, it follows that $e_{N+L} \leq$ $p \leq e_{N}$. We can also estimate $\|\sigma(p)-p\|$ as required (see [15] for details).

If $\tau$ is a tracial state on $A$ as in the statement, then $\tau(f)=\varphi(e)^{2 N+2 L+1}$ and $\tau\left(f_{k}\right)=\tau(f)(1-\varphi(e))$. Since $\tau(v)=0$, we get that $\tau(p)=\tau(f)+\tau(f)(1-\varphi(e)) L=$ $\tau(f)(1+L(1-\varphi(e)))$. This concludes the proof.

Lemma 4.9. For each $i=0,1, \ldots, n$ and $\epsilon>0$ there exists a projection $f_{i} \in A_{n+1}$ such that $\left\|\sigma\left(f_{i}\right)-f_{i}\right\|<\epsilon$ and $\tau_{\lambda}\left(f_{i}\right)$ depends only on $\lambda_{i} \in[0,1]$ and is a decreasing function in $\lambda_{i}$ and for any $\lambda \in \Lambda_{n+1}$,

$$
\begin{array}{ll}
\tau_{\lambda}\left(f_{i}\right)<\epsilon & \text { if } \lambda_{i}>\epsilon \\
\tau_{\lambda}\left(f_{i}\right)=1 & \text { if } \lambda_{i}=0
\end{array}
$$

from which there exists $a \delta>0$ such that if $\lambda_{i}<\delta$ then $\tau_{\lambda}\left(f_{i}\right)>1-\epsilon$. Moreover $f_{i}$ can be chosen from $\left(\otimes_{Z} M_{n+1}\right)^{\gamma}$ for some finite subset $Z$ of $\mathbb{Z}$ and have the following property. If $\lambda \in \Lambda_{n+1}$ satisfies that $\operatorname{dist}\left(\lambda, \operatorname{Ex}\left(\Lambda_{n+1}\right)\right) \geq \epsilon+2 \delta$ and $\operatorname{dist}(\lambda,[i, j])<\delta$ for some $i \neq j$, then $\tau_{\lambda}\left(f_{i}\right)<\epsilon, \tau_{\lambda}\left(f_{j}\right)<\epsilon$, and $\tau_{\lambda}\left(f_{k}\right)>1-\epsilon$ for all $k \neq i, j$.

Proof. Let $E_{i}^{(m)}=\sum_{j \neq i} e_{j j}^{(m)}=1-e_{i i}^{(m)}$. It is shown by the previous lemma that for large $L, N \in \mathbb{N}$ there is a projection $f_{i} \in\left(\bigotimes_{m=-2 L-N}^{2 L+N} M_{n+1}\right)^{\gamma}$ between $\bigotimes_{m=-L-N}^{L+N} E_{i}^{(m)}$ and $\bigotimes_{m=-N}^{N} E_{i}^{(m)}$ such that $\left\|\sigma\left(f_{i}\right)-f_{i}\right\|$ is of order $L^{-1 / 2}$. By using the explicit formula for $f_{i}$ given there, we have that

$$
\tau_{\lambda}\left(f_{i}\right)=\left(1-\lambda_{i}\right)^{2 N+2 L+1}\left(1+L \lambda_{i}\right)
$$

for $\lambda \in \Lambda_{n+1}$. Then if $\lambda_{i}>\epsilon$, then $\tau_{\lambda}(f) \leq \tau_{\lambda}\left(\otimes_{m=-N}^{N} E_{i}^{(m)}\right)<(1-\epsilon)^{2 N+1}$ and if $\lambda_{i}=0$ then $\tau_{\lambda}\left(f_{i}\right)=1$. By choosing $L, N$ sufficiently large, this concludes the proof of the first part.

To show the second part suppose, in contrast, that $\operatorname{dist}\left(\lambda, \operatorname{Ex}\left(\Lambda_{n+1}\right)\right) \geq \epsilon+2 \delta$, $\operatorname{dist}(\lambda,[i, j])<\delta$, and $\tau_{\lambda}\left(f_{i}\right) \geq \epsilon$; from the last condition it follows that $\lambda_{i} \leq \epsilon$. Since $\min _{0 \leq t \leq 1} \max \left\{\left|\lambda_{i}-t\right|,\left|\lambda_{j}-(1-t)\right|\right\}<\delta$, there is a $t \in[0,1]$ such that $t-\lambda_{i}<\delta$ and $1-t-\lambda_{j}<\delta$, which implies that $\lambda_{j}>1-\lambda_{i}-2 \delta \geq 1-\epsilon-2 \delta$. Since this contradicts that $\operatorname{dist}\left(\lambda, \operatorname{Ex}\left(\Lambda_{n+1}\right)\right) \geq \epsilon+2 \delta$, we get that $\tau_{\lambda}\left(f_{i}\right)<\epsilon$. The same argument yields that $\tau_{\lambda}\left(f_{j}\right)<\epsilon$. Since $\lambda_{k}<\delta$ for $k \neq i, j$, it follows that $\tau_{\lambda}\left(f_{k}\right)>1-\epsilon$ for such a $k$.

Continuation of the proof of Lemma 4.6. By the previous lemma for any small $\epsilon>0$ we choose projections $f_{0}, f_{1}, \ldots, f_{n-1}$ in $A_{n+1}$ such that $\tau_{\lambda}\left(f_{i}\right)<\epsilon$ if $\lambda_{i}>\epsilon$, $\tau_{\lambda}\left(f_{i}\right)=1$ if $\lambda_{i}=0$, and $\left\|\sigma\left(f_{i}\right)-f_{i}\right\| \approx 0$ for all $i$. More explicitly we use the same formula (as a function of $E_{i}^{(m)}=1-e_{i i}^{(m)}$ ) for constructing $f_{i}$; hence, we have that $\tau_{\lambda}\left(f_{i}\right)=\tau_{\mu}\left(f_{j}\right)$ if $\lambda_{i}=\mu_{j}$. Furthermore, we assume that all the $f_{i}$ reside at different places (by applying powers of $\sigma$ ), i.e. there are finite subsets $Z_{0}, \ldots, Z_{n-1}$ of $\mathbb{Z}$ which are mutually disjoint such that $f_{i} \in\left(\bigotimes_{Z_{i}} M_{n+1}\right)^{\gamma}$; in particular, the $f_{i}$ commute with each other and $\left[f_{i} f_{j}\right]=\left[f_{i}\right]\left[f_{j}\right]$ in $\mathbb{Z}\left[x_{1}, \ldots, x_{n}\right]$ for $i \neq j$.

We suppose that $\epsilon<(n+2)^{-1}$ and note that $\delta>0$ is chosen so that if $\lambda_{k}<\delta$ then $\tau_{\lambda}\left(f_{k}\right)>1-\epsilon$; obviously $\delta<\epsilon$. Let $N, N_{1}, C, S \in \mathbb{N}$ be as in the proof before the interruption. For each $k=0,1, \ldots, n$ let $\left(E_{i}^{k}\right)$ be an $N_{1}$-cycle in $A_{n+1}$ such that $\left\|\sigma\left(E_{i}^{k}\right)-E_{i+1}^{k}\right\| \approx 0$,

$$
N_{1}\left[E_{0}^{k}\right]\left(x_{1}, \ldots, x_{n}\right)>1-\epsilon
$$

if $\min _{i \neq k} \operatorname{dist}(x,[i, k]) \geq \delta$, and

$$
\left[E_{0}^{k}\right] \geq C\left(x_{0} x_{1} \cdots x_{n}\right)^{S}, \quad\left[E_{0}^{k}\right] \geq c \sum_{i \neq j ; i, j \neq k}\left(x_{i} x_{j}\right)^{S}
$$

for some $c>0$, if $x_{k} \leq 1-n \delta$. We suppose that all $\left(E_{i}^{k}\right)$ and $f_{k}$ reside at disjoint subsets of $\mathbb{Z}$ for various $k$; in particular, we have assumed that $f_{k} E_{i}^{k}$ is a projection and that $\tau_{\lambda}\left(f_{k} E_{i}^{k}\right)=\tau_{\lambda}\left(f_{k}\right) \tau_{\lambda}\left(E_{i}^{k}\right)$. Now we define an $N_{1}$-cycle:

$$
E_{i}=\bigvee_{k=0}^{n}\left(f_{k} E_{i}^{k}\right)+\left(1-\bigvee_{k=0}^{n} f_{k}\right) E_{i}^{0}
$$

Here we have use the fact that all $f_{k} E_{i}^{k}$ commute with each other. We show that $\left(E_{i}\right)$ is the required $N_{1}$-cycle in $A_{n+1}$.

Since $\left\|\sigma\left(f_{k} E_{i}^{k}\right)-f_{k} E_{i+1}^{k}\right\| \approx 0$ and $\|\sigma(F)-F\| \approx 0$ where $F=\bigvee_{k=0}^{n} f_{k}=$ $f_{0}+f_{1}\left(1-f_{1}\right)+f_{2}\left(1-f_{0}-f_{1}\left(1-f_{0}\right)\right)+\cdots$, we obtain that $\left\|\sigma\left(E_{i}\right)-E_{i+1}\right\| \approx 0$.

Let $\lambda \in \Lambda_{n+1}$ with $\lambda_{0}=1-\lambda_{1}-\cdots-\lambda_{n}$. We take $k \in\{0,1, \ldots, n\}$ such that $\lambda_{k}=$ $\min \left\{\lambda_{i} \mid i=0, \ldots, n\right\}$. Then we have that $\tau_{\lambda}\left(f_{k}\right)=\max \left\{\tau_{\lambda}\left(f_{i}\right) \mid i=0, \ldots, n\right\}$ since $\lambda_{i} \mapsto \tau_{\lambda}\left(f_{i}\right)$ is decreasing and $\tau_{\lambda}\left(f_{i}\right)=\tau_{\lambda}\left(f_{j}\right)$ if $\lambda_{i}=\lambda_{j}$. If $\tau_{\lambda}\left(f_{k}\right) \geq 1 /(n+2)>\epsilon$, then we have that $\lambda_{k} \leq \epsilon$.

If $\tau_{\lambda}\left(f_{i}\right)<1 /(n+2)<1-\epsilon$ for all $i$, then $\lambda_{i} \geq \delta$ and $\tau_{\lambda}(F) \leq \sum_{i=0}^{n} \tau_{\lambda}\left(f_{i}\right)<$ $(n+1) /(n+2)<1-1 /(n+2)$. Note that for any $i$,

$$
\lambda_{i}=1-\sum_{j \neq i} \lambda_{j}<1-n \delta .
$$

Thus, we have two cases for $\lambda \in \Lambda_{n+1}$ : (I) the smallest of all $\lambda_{i}$ is given by $\lambda_{k}$ satisfying $\lambda_{k} \leq \epsilon<1-n \delta$ and $\tau_{\lambda}\left(f_{k}\right) \geq 1 /(n+2)$; and (II) $\tau_{\lambda}(1-F)>1 /(n+2), \lambda_{i} \geq \delta$ for all $i$, and $\lambda_{0} \leq 1-n \delta$.

In the first case

$$
\tau_{\lambda}\left(E_{0}\right) \geq \tau_{\lambda}\left(f_{k}\right) \tau_{\lambda}\left(E_{0}^{k}\right) \geq \frac{C}{n+2}\left(\lambda_{0} \lambda_{1} \cdots \lambda_{n}\right)^{S}
$$

In the second case

$$
\tau_{\lambda}\left(E_{0}\right) \geq \tau_{\lambda}(1-F) \tau_{\lambda}\left(E_{0}^{0}\right) \geq \frac{C}{n+2}\left(\lambda_{0} \lambda_{1} \cdots \lambda_{n}\right)^{S}
$$

In either case we have one of the required estimates from below.
To prove the other estimate, we have, for the first case,

$$
\tau_{\lambda}\left(E_{0}\right) \geq \frac{c}{n+2} \sum_{i \neq j ; i, j \neq k}\left(\lambda_{i} \lambda_{j}\right)^{S},
$$

where $\lambda_{k}<1-n \delta$ is used. Since $\lambda_{k}$ is the smallest in $\lambda_{i}$, we have that for each $i \neq k$

$$
\left(\lambda_{i} \lambda_{k}\right)^{S} \leq 1 /(n-1) \sum_{j \neq i, k}\left(\lambda_{i} \lambda_{j}\right)^{S} \leq \sum_{j \neq i, k}\left(\lambda_{i} \lambda_{j}\right)^{S},
$$

which implies that

$$
\begin{aligned}
\sum_{i \neq j}\left(\lambda_{i} \lambda_{j}\right)^{S} & =\sum_{i \neq j ; i, j \neq k}\left(\lambda_{i} \lambda_{j}\right)^{S}+\sum_{i \neq k}\left(\lambda_{i} \lambda_{k}\right)^{S} \\
& \leq \sum_{i \neq j ; i, j \neq k}\left(\lambda_{i} \lambda_{j}\right)^{S}+\sum_{i \neq k} \sum_{j \neq i, k}\left(\lambda_{i} \lambda_{j}\right)^{S} \\
& =3 \sum_{i \neq j ; i, j \neq k}\left(\lambda_{i} \lambda_{j}\right)^{S} .
\end{aligned}
$$

Hence, we have that

$$
\tau_{\lambda}\left(E_{0}\right) \geq \frac{c}{3(n+2)} \sum_{i \neq j}\left(\lambda_{i} \lambda_{j}\right)^{S}
$$

We have, for the second case, that

$$
\tau_{\lambda}\left(E_{0}\right) \geq \frac{c}{n+2} \sum_{i \neq j ; i, j>0}\left(\lambda_{i} \lambda_{j}\right)^{S}
$$

Since $\left(\lambda_{0} \lambda_{i}\right)^{S} \leq \lambda_{i}^{S} \leq(n-1)^{-1} \delta^{-S} \sum_{j \neq i, 0}\left(\lambda_{i} \lambda_{j}\right)^{S} \leq \delta^{-S} \sum_{j \neq i, 0}\left(\lambda_{i} \lambda_{j}\right)^{S}$, we have that

$$
\sum_{i \neq j ; i, j>0}\left(\lambda_{i} \lambda_{j}\right)^{S} \geq \frac{1}{1+2 \delta^{-S}} \sum_{i \neq j}\left(\lambda_{i} \lambda_{j}\right)^{S}
$$

Hence, in this case we get that

$$
\tau_{\lambda}\left(E_{0}\right) \geq \frac{c}{(n+2)\left(1+2 \delta^{-S}\right)} \sum_{i \neq j}\left(\lambda_{i} \lambda_{j}\right)^{S}
$$

Combining these cases we conclude that there is a $c>0$ such that

$$
\left[E_{0}\right]\left(x_{1}, \ldots, x_{n}\right) \geq c \sum_{i \neq j}\left(x_{i} x_{j}\right)^{S}
$$

on $\Lambda_{n+1}$.
Let $\lambda \in \Lambda_{n+1}$ such that $\operatorname{dist}\left(\lambda, \operatorname{Ex}\left(\Lambda_{n+1}\right)\right) \geq 3 \epsilon(>\epsilon+2 \delta)$.

If $\operatorname{dist}(\lambda,[i, j])<\delta$ for some edge $[i, j]$, then by the previous lemma, $\tau_{\lambda}\left(f_{i}\right)<$ $\epsilon, \tau_{\lambda}\left(f_{j}\right)<\epsilon$, and $\tau_{\lambda}\left(f_{k}\right)>1-\epsilon$ for $k \neq i, j$. Hence, for $k \neq i, j$, we have that

$$
N_{1} \tau_{\lambda}\left(E_{0}\right) \geq N_{1} \tau_{\lambda}\left(f_{k}\right) \tau_{\lambda}\left(E_{0}^{k}\right)>(1-\epsilon)^{2}
$$

where we have used the fact that $\operatorname{dist}(\lambda,[k, \ell]) \geq \delta$ for any $\ell$.
If $\operatorname{dist}(\lambda,[i, j]) \geq \delta$ for any edge $[i, j]$, then $\tau_{\lambda}\left(f_{k}\right) \leq 1-\epsilon$ and $\tau_{\lambda}\left(E_{0}^{k}\right)>1-\epsilon$ for all $k$. Since $\tau_{\lambda}\left(1-\bigvee_{k=0}^{n}\left(f_{k} E_{0}^{k}\right)\right)=\prod_{k=0}^{n} \tau_{\lambda}\left(1-f_{k} E_{0}^{k}\right) \leq \prod_{k=0}^{n}\left(1-\tau_{\lambda}\left(f_{k}\right)(1-\epsilon)\right)$, we have that

$$
N_{1} \tau\left(E_{0}\right) \geq 1-\prod_{k=0}^{n}\left(1-\tau_{\lambda}\left(f_{k}\right)(1-\epsilon)\right)+\prod_{k=0}^{n}\left(1-\tau_{\lambda}\left(f_{k}\right)\right) \cdot(1-\epsilon)
$$

If $\epsilon$ is sufficiently small, then the right-hand side is approximately equal to $1-\epsilon \prod_{k=0}^{n}$ $\left(1-\tau_{\lambda}\left(f_{k}\right)\right)\left(1+\sum_{k=0}^{n} \tau_{\lambda}\left(f_{k}\right) /\left(1-\tau_{\lambda}\left(f_{k}\right)\right)\right)$, which is bigger than $1-(n+2) \epsilon$.

This shows that $N_{1}\left[E_{0}\right]\left(x_{1}, \ldots, x_{n}\right)$ is close to 1 uniformly on the compact subset of $\Lambda_{n+1}$ consisting of points distant at least $3 \epsilon$ from the vertices.

The following lemma follows from the proof of the above lemma.
Lemma 4.10. Let $n, m \in \mathbb{N}$ with $n \geq 2$. Suppose that the $n$-shift has the approximate Rohlin property. Then for any $N \in \mathbb{N}$ there are $N_{1}, S \in \mathbb{N}$ and $C>0$ as in Definition 4.2 such that there is a sequence $\left(e_{i}^{(k)}\right)$ of $N_{1}$-cycles in $A_{n+m}$ with $E_{k}=\sum_{i=0}^{N_{1}-1} e_{i}^{(k)}$ and a decreasing sequence $\left(\delta_{k}\right)$ of positive numbers with $\lim _{k} \delta_{k}=0$ satisfying

$$
\begin{gathered}
\max _{i}\left\|\sigma\left(e_{i}^{(k)}\right)-E_{i+1}^{(k)}\right\| \rightarrow 0, \quad k \rightarrow \infty, \\
{\left[E_{k}\right](x) \geq C\left(x_{0} x_{1} \cdots x_{n-1}\right)^{S}, \quad x \in \Lambda_{n+m}^{\left(\delta_{k}\right)},} \\
{\left[E_{k}\right](x) \geq c_{k} \sum_{i \neq j ; i, j<n}\left(x_{i} x_{j}\right)^{S}, \quad x \in \Lambda_{n+m}^{\left(\delta_{k}\right)}}
\end{gathered}
$$

for some $c_{k}>0$, and $\left[E_{k}\right](x)$ converges to 1 uniformly on every compact subset of $\Lambda_{n+m}^{(0)} \backslash \operatorname{Ex}\left(\Lambda_{n+m}\right)$, where

$$
\Lambda_{n+m}^{(\delta)}=\left\{x \in \Lambda_{n+m} \mid \sum_{i=0}^{m-1} x_{n+i}<1-\delta\right\}
$$

for $\delta \geq 0$.
Proof. Let $K, L, M \in \mathbb{N}$ be such that $M \gg L$ and let $Z_{M, K}=\{m M \mid m=$ $0,1, \ldots, K-1\}$ and $Z_{M, K, L}=\bigcup\left\{Z_{M, K}+m \mid m=0,1, \ldots, L\right\}$. We define a (non-unital) embedding $\iota_{0}$ of $M_{n}$ into $\bigotimes_{\mathbb{Z}} M_{n+m}$ by $\iota_{0}\left(e_{i j}\right)=e_{i j}^{(0)}$, where $e_{i j}^{(0)}$ are matrix units of $M_{n+m}$ at $0 \in \mathbb{Z}$. In particular, $f_{0}=1-\iota_{0}(1)=\sum_{i=n}^{n+m-1} e_{i i}^{(0)}$. We then define a homomorphism $\psi_{0}$ into $\bigotimes_{Z_{M, K}} M_{n+m}$ by

$$
\psi_{0}(x)=\iota_{0}(x)+f_{0} \cdot \iota_{M}(x)+f_{0} f_{M} \cdot \iota_{2 M}(x)+\cdots+f_{0} \cdots f_{(K-2) M} \cdot \iota_{(K-1) M}(x),
$$

where $\iota_{M}(x)=\sigma^{M} \iota_{0}(x), f_{M}=\sigma^{M}\left(f_{0}\right)$, etc. By setting $\psi_{k}=\sigma^{k} \psi_{0}$, we define a homomorphism $\psi$ of $\otimes_{0}^{L-1} M_{n}$ into $\otimes_{Z_{M, K, L}} M_{n+m} \subset \otimes_{\mathbb{Z}} M_{n+m}$ by $\psi_{0} \otimes \psi_{1} \otimes \cdots \otimes$ $\psi_{L-1}$.

Then we proceed just as in the proof of Lemma 4.6. In particular, we should note that for a projection $e \in B_{L}=\left(\otimes_{0}^{L-1} M_{n}\right)^{\gamma}$ and $\lambda \in \Lambda_{n+m}$ with $c \equiv \sum_{i=0}^{m-1} \lambda_{n+i}<1$,

$$
\tau_{\lambda} \psi(e)=\left(1-c^{K}\right)^{L} \tau_{\mu}(e),
$$

where $\mu=(1-c)^{-1}\left(\lambda_{1}, \ldots, \lambda_{n-1}\right) \in \Lambda_{n}$.
Lemma 4.11. Let $n, m \in \mathbb{Z}_{+}$and let $S, M \in \mathbb{N}$ with $M>S$. Let e be a non-zero projection in $A_{n+m}$ such that the $i$-vanguard $V^{i}([e])$ for $i=0,1, \ldots, n-1$ is confined in the $S$-deep face generated by the $j$ th vertices for $j=0, \ldots, n-1$, i.e. if $v \in V^{i}([e])$ for $i=0,1, \ldots, n-1$, then $v_{j} \leq S$ for $j<n$ (and $j \neq i$ ) and $v_{j}=0$ for $j \geq n$. Let $f$ be a projection in $I\left(\sum_{i=0}^{n-1} M e_{i}\right)$ such that $[e]-[f]$ is strictly positive on $\Lambda_{n+m}^{(\delta)}=\left\{x \in \Lambda_{n+m} \mid x_{n}+\cdots+x_{n+m-1}<\delta\right\}$ for some $\delta>0$ except for the $n$ vertices of $\Lambda_{n+m-1}$ inside. Then for any $k=\left(k_{0}, k_{1}, \ldots, k_{m-1}\right) \in \mathbb{N}^{m}$ it follows that $[Q(e)] \geq[Q(f)]$, where $Q$ is the quotient map of $A_{n+m}$ onto

$$
B=A_{n+m} / \sum_{i=0}^{m-1} I\left(k_{i} e_{n+i}\right)
$$

Proof. We show that there are $C_{n+i} \in \mathbb{N}$ for $i=0,1, \ldots, m-1$ such that

$$
[e]+\sum_{i=0}^{m-1} C_{n+i}\left(x_{0} x_{1} \cdots x_{n-1}\right)^{M} x_{n+i}^{k_{i}} \geq[f]
$$

in $\mathbb{Z}\left[x_{1}, \ldots, x_{n+m-1}\right]$. Since the order ideal corresponding to $\sum_{i=0}^{m-1} I\left(k_{i} e_{n+i}\right)$ is

$$
\sum_{i=0}^{m-1} x_{n+i}^{k_{i}} \mathbb{Z}\left[x_{1}, \ldots, x_{n+m-1}\right]
$$

this will give the result.
Since $[f]=\left(x_{0} \cdots x_{n-1}\right)^{M} r$ with some $r \in \mathbb{Z}\left[x_{1}, \ldots, x_{n+m-1}\right]$, let $D$ be the maximum of $r$ as a function on $\Lambda_{n+m} \backslash \Lambda_{n+m}^{(\delta)}$. We choose $C_{i} \in \mathbb{N}$ sufficiently large so that

$$
q\left(x_{1}, \ldots, x_{n+m-1}\right)=\sum_{i=0}^{m-1} C_{n+i} x_{n+i}^{k_{i}}>D
$$

on $\Lambda_{n+m} \backslash \Lambda_{n+m}^{(\delta)}$. Note that the $i$-vanguard $V^{i}(q-r)$ of $q-r$ with $i \geq n$ consists of one point 0 . If $q_{1}=\left(x_{0} \cdots x_{n-1}\right)^{M} q$, then it follows that $[e]+q_{1}-[f]$ is strictly positive on $\Lambda_{n+m-1}$ except for the vertices. (If it vanishes at some point, excluding the endpoints, on the edge $[i, j]$ with $i<n, j \geq n$, it vanishes on the whole $[i, j]$, which contradicts the assumption that it is strictly positive near the vertex $x_{i}=1$.) We have to check the positivity condition on the vanguards of $[e]+q_{1}-[f]$ for each $i=0,1, \ldots, n+m-1$.

Let $L \in \mathbb{N}$ be such that $[e]$ and $q_{1}-[f]$ can be expressed in terms of $x^{v}, v \in V_{L}$ :

$$
[e]=\sum_{v \in V_{L}} b_{v} x^{v}, \quad q_{1}-[f]=\sum_{v \in V_{L}} d_{w} x^{w}
$$

where we suppose that $b_{v} \geq 0$ for all $v \in V_{L}$.

Let $v \in V_{L}^{i}\left([e]+q_{1}-[f]\right)$ with $i=0, \ldots, n-1$. If $v_{j}>S$ for some $j \neq i$ in $\{0,1, \ldots, n-1\}$, then $v$ does not belong to the $i$ vanguard $V_{L}^{i}([e])$ of $[e]$ by the assumption. Hence, $v_{j} \geq M$ for all $j \neq i$ in $\{0, \ldots, n-1\}$ (otherwise $d_{w}=0$ for any $w$ in front of $v$ ). Hence, there must be a $v^{\prime} \in V_{L}^{i}([e])$ such that $v^{\prime}$ is in front of $v$ (since $v_{j}^{\prime} \leq S<M \leq v_{j}$ for $j \neq i$ in $\{0,1, \ldots, n-1\}$ and $v_{j}=0$ for $j \geq n$ ). As $b_{v^{\prime}}+d_{v^{\prime}}=b_{v^{\prime}}>0, v$ cannot belong to $V_{L}^{i}\left([e]+q_{1}-[f]\right)$. This shows that if $v \in V_{L}^{i}\left([e]+q_{1}-[f]\right)$, then $v_{j} \leq S$ for all $j \neq i$ in $\{0, \ldots, n-1\}$; so its coefficient is positive because it is non-zero and has contributions only from $[e] ; v$ must belong to $V_{L}^{i}([e])$.

Next let $v \in V_{L}^{i}\left([e]+q_{1}-[f]\right)$ with $i=n, \ldots, n+m-1$. If $v_{j}<M$ for some $j=0, \ldots, n-1$, then $d_{v}=0$ and hence $b_{v}+d_{v}=b_{v}$ must be positive. If $v_{j} \geq M$ for all $j=0,1, \ldots, n-1$, then $v$ must be $\sum_{j=0}^{n-1} M e_{j}+(L-n M) e_{i}$, where $d_{v}>0$ (which is the value taken by $q-r$ on $x_{i}=1$ ); hence $b_{v}+d_{v} \geq d_{v}>0$.

Thus, we conclude the proof that the coefficients of $[e]+q_{1}-[f]$ takes positive values on the vanguards. Hence, $[e]+q_{1}-[f] \geq 0$ as claimed.

The following lemma, together with Lemmas 4.3, 4.4 and 4.6, completes the induction, thus implying the main theorem of this note. We closely follow the arguments given in the proof of Theorem 2.8 of [19], where we must apologize since there was some confusion in details.
Lemma 4.12. Let $n \in \mathbb{N}$ with $n \geq 2$. If the $n$-shift has the approximate Rohlin property and $(n-1)$-shift has the stable one-cocycle property, then the $n$-shift has the stable one-cocycle property.

Proof. Let $m \in \mathbb{Z}_{+}$and $k=\left(k_{0}, k_{1}, \ldots, k_{m-1}\right) \in \mathbb{N}^{m}$. We show that the shift $\sigma$ on the quotient

$$
B=A_{n+m} / \sum_{i=0}^{m-1} I\left(k_{i} e_{n+i}\right)
$$

has the one-cocycle property.
Let $u \in \mathcal{U}(B)$ be such that $\varphi_{i}(u)=1$ for $i=0,1, \ldots, n-1$, where $\varphi_{i}$ is the character on $B$ induced from the character on $A_{n}$ denoted by the same symbol through

$$
A_{n} \cong A_{n+m} / \sum_{i=0}^{m-1} I\left(e_{n+i}\right)
$$

We have to show that for any $\epsilon>0$ there is a $v \in \mathcal{U}(B)$ such that $\left\|u-v \sigma(v)^{*}\right\|<\epsilon$ and also the two other versions. First we just concentrate on approximating $u$ by $v \sigma(v)^{*}$.

We denote by $Q$ the quotient map of $A_{n+m}$ onto $B$. Let $M \in \mathbb{N}$. From the short exact sequence

$$
0 \rightarrow Q\left(I\left(M e_{0}\right)\right) \rightarrow B \rightarrow B / Q\left(I\left(M e_{0}\right)\right) \rightarrow 0
$$

and the assumption that the shift has the one-cocycle property for $B / Q\left(I\left(M e_{0}\right)\right) \cong$ $A_{n+m} /\left(I\left(M e_{0}\right)+\sum_{i=0}^{m-1} I\left(k_{i} e_{n+i}\right)\right)$, there is a $v \in \mathcal{U}(B)$ such that

$$
\left\|u-v \sigma\left(v^{*}\right)+Q\left(I\left(M e_{0}\right)\right)\right\| \approx 0
$$

where we have used that $\varphi_{i}(Q(u))=1$ for $i=1, \ldots, n-1$. Thus, by taking a unitary in $Q\left(I\left(M e_{0}\right)\right)+1$ close to $v^{*} u \sigma(v)$ instead of $u$, we may suppose that $u \in Q\left(I\left(M e_{0}\right)\right)+1$.

Similarly, from the short exact sequence

$$
0 \rightarrow Q\left(I\left(M e_{0}+M e_{1}\right)\right) \rightarrow Q\left(I\left(M e_{0}\right)\right) \rightarrow Q\left(I\left(M e_{0}\right)\right) / Q\left(I\left(M e_{0}+M e_{1}\right)\right) \rightarrow 0
$$

where the quotient is isomorphic to $I\left(M e_{0}\right) /\left(I\left(M e_{0}+M e_{1}\right)+\sum_{i=0}^{m-1} I\left(M e_{0}+k_{i} e_{n+i}\right)\right)$ which is an ideal of $A_{n+m} /\left(I\left(M e_{1}\right)+\sum_{i=0}^{m-1} I\left(k_{i} e_{n+i}\right)\right)$, there is a unitary $v \in$ $Q\left(I\left(M e_{0}\right)\right)+1$ such that

$$
\left\|u-v \sigma\left(v^{*}\right)+Q\left(I\left(M e_{0}+M e_{1}\right)\right)\right\| \approx 0
$$

Thus, we may assume that $u \in Q\left(I\left(M e_{0}+M e_{1}\right)\right)+1$. Repeating this process, we reach the following conclusion. For any unitary $u \in \mathcal{U}(B), M \in \mathbb{N}$, and $\epsilon>0$ we find a unitary $u^{\prime} \in Q\left(J_{M}\right)+1$ and a unitary $v \in B$ such that $\left\|v^{*} u \sigma(v)-u^{\prime}\right\|<\epsilon$, where $J_{M}=I\left(\sum_{i=0}^{n-1} M e_{i}\right)$. We may further suppose that there is a projection $f \in Q\left(J_{M}\right)$ such that $\sigma(f) \approx f$ and $u^{\prime}=u^{\prime} f+1-f$.

We now have to approximate $u^{\prime}$ by $v \sigma(v)^{*}$ for some $v \in \mathcal{U}(B)$ by assuming that $M$ is sufficiently large.

We proceed as follows. By mapping a Rohlin cycle for the shift on $A_{n}$ into $A_{n+m}$ (by Lemma 4.10) and then into the quotient $B$, we try to approximate $u^{\prime}$ by $v \sigma(v)^{*}$. The Rohlin cycle, so-embedded, must commute with $u^{\prime}$ and some translates of $u^{\prime}$ under $\sigma$. If the sum $E$ of projections in the Rohlin cycle dominated $f$, the support of $u^{\prime}$, then we would be finished by using the now-standard arguments based on the Rohlin cycles. However, $E$ is not the identity and will never dominate $f$, because to assure the commutativity we have to make the Rohlin cycle reside outside of where $u^{\prime}$ resides. Thus, the problem is how to deal with the left-over part $u^{\prime} f(1-E)$ or $u^{\prime \prime}=u^{\prime} f(1-E)+1-$ $f(1-E)$.

To approximate $u^{\prime \prime}$ by $v \sigma(v)^{*}$ we use another similar method. We construct Rohlin cycles for $\sigma$ and for $\operatorname{Ad} u^{\prime \prime} \circ \sigma$ such that they are of the same type and the sum of projections covers $f(1-E)$. For example, let $\left(e_{i}\right)$ and $\left(e_{i}^{\prime}\right)$ be $N$-cycles such that $F=\sum_{i=0}^{N-1} e_{i}=\sum_{i=0}^{N-1} e_{i}^{\prime} \geq f(1-E), \sigma\left(e_{i}\right) \approx e_{i+1}$, and $\operatorname{Ad} u^{\prime \prime} \sigma\left(e_{i}^{\prime}\right) \approx e_{i+1}^{\prime}$; then take a partial isometry $w$ such that $w^{*} w=e_{0}$ and $w w^{*}=e_{0}^{\prime}$, and the desired unitary $v$ can be obtained by modifying

$$
\sum_{i=0}^{N-1}\left(u^{\prime \prime} \sigma\right)^{i}(v)+1-F
$$

(As a matter of fact we need two Rohlin cycles for each of $\sigma$ and $\operatorname{Ad} u^{\prime \prime} \sigma$ to make $F$ large enough to cover $f(1-E)$, and we apply this method to $u^{\prime}$ directly.)

We now denote by $u$ the unitary in $Q\left(J_{M}\right)+1$ for some $M \in \mathbb{N}$ which we have to approximate by a unitary of the form $v \sigma(v)^{*}$.

Let $\epsilon>0$ and $\delta>0$; we specify $\delta$ for the given $\epsilon$ later.
First we choose $N_{1} \in \mathbb{N}$ such that $2 \pi N_{1}^{-1}<\epsilon$; we construct a set of Rohlin cycles of length $N_{1}, N_{1}+1$, and of longer length for each of $\sigma$ and $\operatorname{Ad} u^{\prime} \sigma$ to approximate $u^{\prime}$ by a unitary of the form $v \sigma(v)^{*}$ within the error of order $\epsilon$.

We then choose $N_{2} \in \mathbb{N}$ such that if $U \in M_{N_{2}+1}$ is a unitary with eigenvalues $\exp \left(2 \pi i k / N_{2}\right) ; k=0,1, \ldots, N_{2}$, then there is a $N_{1}$-cycle $\left(e_{i}\right)$ and a $N_{1}+1 \operatorname{cycle}\left(e_{i}^{\prime}\right)$
in $M_{N_{2}+1}$ such that $\sum_{i=0}^{N_{1}-1} e_{i}+\sum_{i=0}^{N_{1}} e_{i}^{\prime}=1$ and

$$
\left\|\operatorname{Ad} U\left(e_{i}\right)-e_{i+1}\right\|<\delta, \quad\left\|\operatorname{Ad} U\left(e_{i}^{\prime}\right)-e_{i+1}^{\prime}\right\|<\delta .
$$

By making $N_{2}$ larger if necessary, we obtain a sequence $\left(e_{2 i}^{(k)}\right)$ of $N_{2}$-cycles in $A_{n+m}$ with $E_{2, k}=\sum_{i} e_{2, i}^{(k)}$ such that

$$
\begin{gathered}
\max _{i}\left\|\sigma\left(e_{2, i}^{(k)}\right)-e_{2, i+1}^{(k)}\right\|<\delta_{k}, \\
{\left[E_{2, k}\right](x) \geq C\left(x_{0} x_{1} \cdots x_{n-1}\right)^{M_{2}}, \quad x \in \Lambda_{n+m}^{\left(\delta_{k}\right)},} \\
{\left[E_{2, k}\right](x) \geq c_{k} \sum_{i \neq j ; i, j<n}\left(x_{i} x_{j}\right)^{M_{2}}, \quad x \in \Lambda_{n+m}^{\left(\delta_{k}\right)},}
\end{gathered}
$$

and $\left[E_{2, k}\right](x)$ converges to 1 uniformly on every compact subset of $\Lambda_{n+m}^{(0)} \backslash \operatorname{Ex}\left(\Lambda_{n+m}\right)$, where $C>0$ is a constant independent of $k$ as well as $M_{2} \in \mathbb{N}, c_{k}>0$ is a constant, $\left(\delta_{k}\right)$ is a sequence of positive numbers decreasing to 0 , and $\Lambda_{n+m}^{(\delta)}=\left\{x \in \Lambda_{n+m}\right\}$ $\left.\sum_{i=0}^{m-1} x_{n+i}<\delta\right\}$.

Let $\delta^{\prime}>0$ be sufficiently small and let $N_{3} \in \mathbb{N}$ be such that $2 N_{3}^{-1 / 2}<\delta^{\prime} ; \delta^{\prime}$ will be chosen for $\delta$ and $N_{2}$. By taking $N^{\prime}$-cycles for $\sigma$ for $N^{\prime} \geq N_{2} N_{3}$, we obtain a sequence $\left(e_{3 i}^{(k)}\right)$ of $N_{3}$ projections in $A_{n+m}$ with $E_{3, k}=\sum_{i} e_{3, i}^{(k)}$ such that

$$
\begin{gathered}
\max _{i}\left\|\sigma^{N_{2}}\left(e_{3, i}^{(k)}\right)-e_{3, i+1}^{(k)}\right\|<\delta_{k}, \quad i=0, \ldots, N_{3}-2, \\
{\left[E_{3, k}\right](x) \geq C^{\prime}\left(x_{0} x_{1} \cdots x_{n-1}\right)^{M_{3}}, \quad x \in \Lambda_{n+m}^{\left(\delta_{k}\right)},} \\
{\left[E_{3, k}\right](x) \geq c_{k}^{\prime} \sum_{i \neq j ; i, j<n}\left(x_{i} x_{j}\right)^{M_{3}}, \quad x \in \Lambda_{n+m}^{\left(\delta_{k}\right)},}
\end{gathered}
$$

and $\left[E_{3, k}\right](x)$ becomes larger than $1-1 / N_{3}$ uniformly on every compact subset of $\Lambda_{n+m}^{(0)} \backslash \operatorname{Ex}\left(\Lambda_{n+m}\right)$, where $C^{\prime}>0$ is a constant independent of $k$ as well as $M_{3} \in \mathbb{N}$, and $c_{k}^{\prime}>0$ is a constant. By taking a smaller one, we set $C^{\prime}=C$ and $c_{k}=c_{k}^{\prime}$.

We set $M=M_{2}+M_{3}+1$. We now assume that the unitary $u \in B$ belongs to $Q\left(J_{M}\right)+1$. We also assume that there is a projection $f \in J_{M}$ such that $\|\sigma(f)-f\|<N_{2}^{-1} \delta^{\prime}$ and $u=u Q(f)+1-Q(f)$ and $f$ is local in the sense that $f$ belongs to $A_{n+m}\left(Z_{1}\right)=$ $\left(\otimes_{Z_{1}} M_{n+m}\right)^{\gamma}$ for some finite subset $Z_{1}$. Further we may assume that $\left(e_{2, i}^{(k)}\right)_{i}$ and $\left(e_{3, i}^{(k)}\right)_{i}$ are local, and they as well as $f$ reside at disjoint subsets.

Then we have that $\left[f\left(1-E_{2, k}\right)\right]=[f]\left[1-E_{2, k}\right],[f] \in\left(x_{0} x_{1} \cdots x_{n-1}\right)^{M} \mathbb{Z}\left[x_{1}, \ldots\right.$, $x_{n+m-1}$ ],

$$
\left[f\left(1-E_{2, k}\right)\right](x) \leq D\left(x_{0} x_{1} \cdots x_{n-1}\right)^{M}, \quad x \in \Lambda_{n+m}
$$

for some $D>0$ and $\left[f\left(1-E_{2, k}\right)\right](x)$ converges to 0 uniformly on every compact subset of $\Lambda_{n+m}^{(0)} \backslash \operatorname{Ex}\left(\Lambda_{n+m}\right)$ and hence of $\Lambda_{n+m}^{(0)}$. On the other hand, we have that

$$
\begin{array}{ll}
{\left[e_{20}^{(k)} e_{30}^{(k)}\right] \geq C^{2}\left(x_{0} x_{1} \cdots x_{n-1}\right)^{M-1},} & x \in \Lambda_{n+m}^{\left(\delta_{k}\right)}, \\
{\left[e_{20}^{(k)} e_{30}^{(k)}\right] \geq c_{k}^{2} \sum_{i \neq j ; i, j<n}\left(x_{i} x_{j}\right)^{M-1},} & x \in \Lambda_{n+m}^{\left(\delta_{k}\right)},
\end{array}
$$

and $\left[e_{20}^{(k)} e_{30}^{(k)}\right](x)$ becomes larger than $\left(N_{2} N_{3}\right)^{-1}\left(1-1 / N_{3}\right)$ uniformly on every compact subset of $\Lambda_{n+m}^{(0)} \backslash \operatorname{Ex}\left(\Lambda_{n+m}\right)$.

If $e_{2, i}^{(k)}$ is constructed from a projection $e_{i}^{\prime} \in A_{n}$ with $\left[e_{i}^{\prime}\right]=g \in \mathbb{Z}\left[x_{1}, \ldots, x_{n-1}\right]$ through the mapping $\phi$ into $A_{n+m}$ (see the proof of Lemma 4.6), then $\left[e_{2, i}^{(k)}\right]$ is given by

$$
\left(1-X^{K}\right)^{L^{\prime}}\left(1+L^{\prime} X\right) \cdot\left(1-X^{K}\right)^{L} g\left(x_{1}(1-X)^{-1}, \ldots, x_{n-1}(1-X)^{-1}\right)
$$

where $L^{\prime}$ depends only on the order of $\left\|\sigma \phi\left(e_{i}^{\prime}\right)-\phi \sigma\left(e_{i}^{\prime}\right)\right\|$ and

$$
X=\sum_{i=n}^{n+m-1} x_{i}
$$

Thus, if $g(x)=g\left(x_{1}, \ldots, x_{n-1}\right)=\sum_{v \in V_{L}} c_{v} x^{v},\left[e_{20}^{(k)}\right]$ is given as

$$
(1+X p(X)) G_{L}\left(x_{1}, \ldots, x_{n+m-1}\right)
$$

where $p(X)$ is a polynomial in $X=x_{n}+\cdots+x_{n+m-1}$ and

$$
G_{L}\left(x_{1}, \ldots, x_{n+m-1}\right)=\sum_{v \in V_{L}}\left(1-x_{1}-\cdots-x_{n+m-1}\right)^{v_{0}} x_{1}^{v_{1}} \cdots x_{n-1}^{v_{n-1}}
$$

Hence, for each $i=0, \ldots, n-1$, the $i$-vanguard $V^{i}\left(\left[e_{20}^{(k)}\right]\right)$ of $\left[e_{20}^{(k)}\right]$ is the same as that of $G_{L}$. More precisely $V^{i}\left(\left[e_{20}^{(k)}\right]\right)$ is confined in the $M_{2}$-deep face generated by the $i$ th vertices for $i=0, \ldots, n-1$, i.e. in the set of $v:\{0,1, \ldots, n+m-1\} \backslash\{i\} \rightarrow \mathbb{Z}_{+}$with $v_{j} \leq M_{2}$ for $j<n($ and $j \neq i)$ and $v_{j}=0$ for $j \geq n$.

Since

$$
\left\{x \in \Lambda_{n+m} \left\lvert\, \sum_{i=0}^{m-1} x_{n+i}<\frac{1}{2}\right., D\left(x_{0} \cdots x_{n-1}\right)^{M}<C^{2}\left(x_{0} \cdots x_{n-1}\right)^{M-1}\right\}
$$

contains a small neighborhood of each vertex $x_{i}=1$ (except for the vertex) for $i=$ $0,1, \ldots, n-1$, we have that for all sufficiently large $k$,

$$
\left[f\left(1-E_{2, k}\right)\right](x)<\left[e_{2,0}^{(k)} e_{3,0}^{(k)}\right](x)
$$

is strict on $\left\{x \in \Lambda_{n+m} \left\lvert\, \sum_{i=0}^{m-1} x_{n+i}<\frac{1}{2}\right.\right\}$ except for the $n$ vertices $x_{i}=1$ for $i=0, \ldots, n-1$. Note also the $i$-vanguards of $\left[e_{2,0}^{(k)} e_{3,0}^{(k)}\right]$ with $i=0, \ldots, n-1$ are confined in the $M_{2}+M_{3}$ deep face generated by the $i$ th vertices for $i=0, \ldots, n-1$. Hence, we can conclude by Lemma 4.11 that

$$
\left[Q\left(f\left(1-E_{2, k}\right)\right)\right] \leq\left[Q\left(e_{2,0}^{(k)} e_{3,0}^{(k)}\right)\right]
$$

for all sufficiently large $k$.
The rest of the arguments proceed as in the proof of Theorem 2.8 of [19].
From now on we will work in $B$ so that we denote $Q(f), Q\left(e_{2, i}^{(k)}\right)$, etc., simply by $f, e_{2, i}^{(k)}$ etc. There should be no confusion.

We take a sufficiently large $k$. First we find a partial isometry onto $f\left(1-E_{2, k}\right) \in B$ from a subprojection of $e_{2,0}^{(k)} \in B$, which is almost $\sigma^{N_{2}}$-invariant, by using $\left(e_{3, i}^{(k)}\right)$ in $B$ as follows.

Let $b$ be a partial isometry in $B=A_{n+m} / \sum_{i=0}^{m-1} I\left(k_{i} e_{n+i}\right)$ such that

$$
b b^{*}=f\left(1-E_{2, k}\right), \quad b^{*} b \leq e_{2,0}^{(k)} e_{3,0}^{(k)} .
$$

Note that we have assumed that $f,\left(e_{2, i}^{(k)}\right)$ and $\left(e_{3, i}^{(k)}\right)$ reside at disjoint subsets of $\mathbb{Z}$, i.e. there are finite subsets $Z_{1}, Z_{2}, Z_{3}$ of $\mathbb{Z}$, which are mutually disjoint, such that $f \in B\left(Z_{1}\right)=$ $Q\left(A_{n+m}\left(Z_{1}\right)\right)=Q\left(\otimes_{Z_{1}} M_{n+m}\right)^{\gamma}, e_{2, i}^{(k)} \in B\left(Z_{2}\right)$, and $e_{3, i}^{(k)} \in B\left(Z_{3}\right)$. (Here we do not mean that $Z_{1}, Z_{2}, Z_{3}$ are independent of $k$; we pick up one $k$ eventually.) We may suppose that there are $x, y, z \in \mathcal{U}(B)$ such that $x \in B\left(Z_{1}\right), y \in B\left(Z_{2}\right), z \in B\left(Z_{3}\right)$, $\|x-1\|$ is of order $\|\sigma(f)-f\|<N_{2}^{-1} \delta^{\prime},\|y-1\| \approx 0,\|z-1\| \approx 0, \operatorname{Ad} x \sigma(f)=f$, $\operatorname{Ad} y \sigma\left(e_{2, i}^{(k)}\right)=e_{2, i+1}^{(k)}$ and $\operatorname{Ad} z(\operatorname{Ad} y \sigma)^{N_{2}}\left(e_{3, i}^{(k)}\right)=\operatorname{Ad} z \sigma^{N_{2}}\left(e_{3, i}^{(k)}\right)=e_{3, i+1}^{(k)}$ for $i<N_{3}-1$. (If the original $Z_{1}, Z_{2}, Z_{3}$ are not only disjoint, but mutually far away by $N_{2}$, this is certainly possible.) Note that the projection $f \in Q\left(J_{M}\right)$ is chosen so that $u=u f+1-f$. We define

$$
d=N_{3}^{-1 / 2} \sum_{i=0}^{N_{3}-1}\left(\operatorname{Ad} z(\operatorname{Ad}(x y) \sigma)^{N_{2}}\right)^{i}(b) .
$$

Then it follows that $d$ is indeed a partial isometry such that $d d^{*}=f\left(1-E_{2, k}\right), d^{*} d \leq e_{2,0}^{(k)}$ and $\left\|\operatorname{Ad} z(\operatorname{Ad}(x y) \sigma)^{N_{2}}(d)-d\right\|<2 N_{3}^{-1 / 2}<\delta^{\prime}$. By assuming that $k$ is sufficiently large, i.e. $\|y-1\|$ and $\|z-1\|$ are sufficiently small, and by the assumption that $\|x-1\| \lesssim N_{2}^{-1} \delta^{\prime}$, we can assume that the unitary $w \in B$ defined through $(\operatorname{Ad}(x y) \sigma)^{N_{2}}=\operatorname{Ad} w \sigma^{N_{2}}$ satisfies $\|w-1\| \lesssim \delta^{\prime}$. Since $\operatorname{Ad} w \sigma^{N_{2}}\left(e_{2,0}^{(k)}\right)=e_{2,0}^{(k)}$ and $\left\|\operatorname{Ad} w \sigma^{N_{2}}\left(d^{*} d\right)-d^{*} d\right\| \lesssim 2 \delta^{\prime}$, there is a unitary $\zeta \in B$ such that $\zeta=\zeta e_{2,0}^{(k)}+1-e_{2,0}^{(k)}$ and $\|\zeta-1\| \lesssim \delta^{\prime}$ such that $\operatorname{Ad} \zeta \operatorname{Ad} w \sigma^{N_{2}}\left(d^{*} d\right)=d^{*} d$. In any case we have that $\left\|\sigma^{N_{2}}(d)-d\right\|$ is at most of order of $\delta^{\prime}$.

By using the partial isometry $d$, we define $D$ to be the $C^{*}$-algebra $D$ generated by $(\operatorname{Ad}(\zeta x y) \sigma)^{k}(d), k=0,1, \ldots, N_{2}-1$. Then $D$ is isomorphic to $M_{N_{2}+1}$ and its identity is left invariant under $\operatorname{Ad}(\zeta x y) \sigma$, since $\operatorname{Ad}(\zeta x y) \sigma\left(d d^{*}\right)=d d^{*}=f\left(1-E_{2, k}\right)$ and $(\operatorname{Ad}(\zeta x y) \sigma)^{N_{2}}\left(d^{*} d\right)=d^{*} d$. Since $\left\|\zeta \sigma^{\prime}(\zeta)\left(\sigma^{\prime}\right)^{2}(\zeta) \cdots\left(\sigma^{\prime}\right)^{N_{2}-1}(\zeta)-1\right\|$ is of order of $\delta^{\prime}$, where $\sigma^{\prime}=\operatorname{Ad}(x y) \sigma$, we have that $\left\|(\operatorname{Ad}(\zeta x y) \sigma)^{N_{2}}(d)-d\right\|$ is of order of $\delta^{\prime}$. Note that if $(\operatorname{Ad}(\zeta x y) \sigma)^{N_{2}}(d)=d$ were true, then $\operatorname{Ad}(\zeta x y) \sigma$ would leave $D$ invariant and be implemented by a unitary $U$ with eigenvalues $\exp \left(2 \pi k N_{2}^{-1}\right), k=0,1, \ldots, N_{2}$. Since $\|\zeta x y-1\|$ is of the order of $\delta^{\prime}$, it follows that $\|(\sigma-\operatorname{Ad} U) \mid D\|<\delta$ with some unitary $U$ as above for a suitable choice of $\delta^{\prime}$ (which can depend on $\delta$ and $N_{2}$ ). Then by the choice of $N_{2}$ we obtain a $N_{1}$-cycle $\left(f_{1, i}\right)$ and a $\left(N_{1}+1\right)$-cycle $\left(f_{2, i}\right)$ in $D$ such that $\sum_{i} f_{1, i}+\sum_{i} f_{2, i}=1_{D}$ and

$$
\left\|\sigma\left(f_{1, i}\right)-f_{1, i+1}\right\|<2 \delta, \quad\left\|\sigma\left(f_{2, i}\right)-f_{2, i+1}\right\|<2 \delta
$$

Let $f_{3, i}=e_{2, i}^{(k)}-(\operatorname{Ad}(\zeta x y) \sigma)^{i}\left(d^{*} d\right)$ for $i=0,1, \ldots, N_{2}-1$. Then it follows that

$$
\operatorname{Ad}(\zeta x y) \sigma\left(f_{3, i}\right)=f_{3, i+1}
$$

hence ( $f_{3, i}$ ) forms a $N_{2}$-cycle in $B$ such that

$$
\left\|\sigma\left(f_{3, i}\right)-f_{3, i+1}\right\|<\delta+\delta^{\prime}<2 \delta .
$$

Thus, we have obtained the three cycles $\left(f_{j, i}\right), j=1,2,3$, which in particular satisfy that

$$
\sum_{i=0}^{N_{1}-1} f_{1, i}+\sum_{i=0}^{N_{1}} f_{2, i}+\sum_{i=0}^{N_{2}-1} f_{3, i}=f \vee E_{2, k} .
$$

We apply the same argument to $\operatorname{Ad} u \sigma$ instead of $\sigma$. Note that $\operatorname{Ad}(x y u) \sigma\left(e_{2, i}^{(k)}\right)=e_{2, i+1}^{(k)}$, $\operatorname{Ad}(x y u) \sigma\left(f\left(1-E_{2, k}\right)\right)=f\left(1-E_{2, k}\right)$ and $\operatorname{Ad} z(\operatorname{Ad}(x y u) \sigma)^{N_{2}}\left(e_{3, i}^{(k)}\right)=\operatorname{Ad} z \sigma^{N_{2}}\left(e_{3, i}^{(k)}\right)$ $=e_{3, i+1}^{(k)}$. With the same $b$ as above, we define a partial isometry $d^{\prime}$ onto $f\left(1-E_{2, k}\right)$ from a subprojection of $e_{2,0}^{(k)}$ by

$$
d^{\prime}=N_{3}^{-1 / 2} \sum_{i=0}^{N_{3}-1}\left(\operatorname{Ad} z(\operatorname{Ad}(x y u) \sigma)^{N_{2}}\right)^{i}(b)
$$

which satisfies that $\left\|\operatorname{Ad} z(\operatorname{Ad}(x y u) \sigma)^{N_{2}}\left(d^{\prime}\right)-d^{\prime}\right\|<2 N_{3}^{-1}<\delta^{\prime}$, where we should note that $\left\|\operatorname{Ad} z(\operatorname{Ad}(x y u) \sigma)^{N_{2}}-\sigma^{N_{2}}\right\| \lesssim \delta^{\prime}$. Let $\zeta^{\prime} \in B$ be a unitary such that $\zeta^{\prime}=\zeta^{\prime} e_{2,0}^{(k)}+1-e_{2,0}^{(k)}$ and $\operatorname{Ad} \zeta^{\prime}(\operatorname{Ad}(x y u) \sigma)^{N_{2}}\left(\left(d^{\prime}\right)^{*} d^{\prime}\right)=\left(d^{\prime}\right)^{*} d^{\prime}$ and $\left\|\zeta^{\prime}-1\right\|$ is of order of $\delta^{\prime}$. Let $D^{\prime}$ denote the $C^{*}$-algebra generated by $\left(\operatorname{Ad}\left(\zeta^{\prime} x y u\right) \sigma\right)^{i}\left(d^{\prime}\right), i=0,1, \ldots, N_{2}-1$. Note that $\left[1_{D}\right]=\left(N_{2}+1\right)\left[f\left(1-E_{2, k}\right)\right]=\left[1_{D^{\prime}}\right]$ in $K_{0}(B)$. Hence, by using the same formula used to define $f_{1, i}$ and $f_{2, i}$ we obtain a $N_{1}$-cycle $\left(f_{1, i}^{\prime}\right)$ and a $\left(N_{1}+1\right)$-cycle $\left(f_{2, i}^{\prime}\right)$ in $D^{\prime}$ such that $\left[f_{1, i}^{\prime}\right]=\left[f_{1, i}\right],\left[f_{2, i}^{\prime}\right]=\left[f_{2, i}\right], \sum_{i} f_{1, i}^{\prime}+\sum_{i} f_{2, i}^{\prime}=1_{D^{\prime}}$, and

$$
\left\|\operatorname{Ad} u \sigma\left(f_{1, i}^{\prime}\right)-f_{1, i+1}^{\prime}\right\|<2 \delta, \quad\left\|\operatorname{Ad} u \sigma\left(f_{2, i}^{\prime}\right)-f_{2, i+1}^{\prime}\right\|<2 \delta .
$$

Let $f_{3, i}^{\prime}=e_{2, i}^{(k)}-\left(\operatorname{Ad}\left(\zeta^{\prime} x y u\right) \sigma\right)^{i}\left(\left(d^{\prime}\right)^{*} d^{\prime}\right)$, whence $\operatorname{Ad}\left(\zeta^{\prime} x y u\right) \sigma\left(f_{3, i}^{\prime}\right)=f_{3, i+1}^{\prime}$. Then $\left(f_{3, i}\right)$ forms a $N_{2}$-cycle such that

$$
\left\|\sigma\left(f_{3, i}\right)-f_{3, i+1}\right\|<2 \delta
$$

Hence, we get the three cycles $\left(f_{j, i}^{\prime}\right), j=1,2,3$, which in particular satisfy that $\left[f_{j, i}^{\prime}\right]=$ [ $f_{j, i}$ ] and that the sum of all the projections $f_{j, i}^{\prime}$ is $f \vee E_{2, k}$.

After having these cycles with appropriate permutation property for $\sigma$ and $\operatorname{Ad} u \sigma$, we proceed exactly as in [19]. We choose partial isometries $b_{1}, b_{2}, b_{3} \in B$ such that

$$
b_{j} b_{j}^{*}=e_{j, 0}^{\prime}, \quad b_{j}^{*} b_{j}=e_{j, 0} .
$$

Then the unitary $v$ which satisfies that $u \sigma(v) \approx v$ is obtained by modifying

$$
v_{1}=\sum_{i=0}^{N_{1}-1}\left(L_{u} \sigma\right)^{i}\left(b_{1}\right)+\sum_{i=0}^{N_{1}}\left(L_{u} \sigma\right)^{i}\left(b_{2}\right)+\sum_{i=0}^{N_{2}-1}\left(L_{u} \sigma\right)^{i}\left(b_{3}\right),
$$

where $L_{u}$ denotes the left multiplication of $u$. The necessary modifications are done as follows. First we choose $Y, Y^{\prime} \in \mathcal{U}(B)$ such that $\|Y-1\| \approx 0,\left\|Y^{\prime}-1\right\| \approx 0$, $\operatorname{Ad} Y \sigma\left(f_{j, i}\right)=f_{j, i+1}$ and $\operatorname{Ad}\left(Y^{\prime} u\right) \sigma\left(f_{j, i}^{\prime}\right)=f_{j, i+1}^{\prime}$. Such $Y$ (respectively $Y^{\prime}$ ) can be obtained by modifying $\zeta x y$ (respectively $\zeta^{\prime} x y$ ). Second we choose sequences $\left(z_{1 i}\right)_{i=0}^{N_{1}-1}$ in $\mathcal{U}\left(f_{1,0} B f_{1,0}\right),\left(z_{2 i}\right)_{i=0}^{N_{1}}$ in $\mathcal{U}\left(f_{2,0} B f_{2,0}\right)$ and $\left(z_{3 i}\right)_{i=0}^{N_{2}-1}$ in $f_{3,0} B f_{3,0}$ such that $z_{1, N_{1}-1}=$ $f_{1,0}, z_{2, N_{1}}=f_{2,0}, z_{3, N_{2}-1}=f_{3,0}$,

$$
\begin{aligned}
& z_{1,0}=b_{1}^{*}\left(L_{Y^{\prime} u} R_{Y^{*}} \sigma\right)^{N_{1}}\left(b_{1}\right), \\
& z_{2,0}=b_{2}^{*}\left(L_{Y^{\prime} u} R_{Y^{*}} \sigma\right)^{N_{1}+1}\left(b_{2}^{*}\right), \\
& z_{3,0}=b_{3}^{*}\left(L_{Y^{\prime} u} R_{Y^{*} \sigma} \sigma\right)^{N_{2}}\left(b_{3}\right),
\end{aligned}
$$

and

$$
\left\|z_{j, i}-z_{j, i+1}\right\|<2 \pi N_{1}^{-1}<\epsilon
$$

where $L_{a}$ (respectively $R_{a}$ ) denotes the left (respectively right) multiplication of $a$. We then define a unitary $v$ by

$$
\begin{aligned}
v= & \sum_{i=0}^{N_{1}-1}\left(L_{Y^{\prime} u} R_{Y^{*}} \sigma\right)^{i}\left(b_{1} z_{1 i}\right)+\sum_{i=0}^{N_{1}}\left(L_{Y^{\prime} u} R_{Y *} \sigma\right)^{i}\left(b_{2} z_{2 i}\right) \\
& +\sum_{i=0}^{N_{2}-1}\left(L_{Y^{\prime} u} R_{Y^{*}} \sigma\right)^{i}\left(b_{3} z_{3 i}\right)+1-f \vee E_{2, k} .
\end{aligned}
$$

Then it follows that $v$ is almost invariant under $L_{Y^{\prime} u} R_{Y^{*}} \sigma$ (up to the order of $\epsilon>2 \pi / N_{1}$ ). Since $Y \approx 1$ and $Y^{\prime} \approx 1$, it follows that $u \sigma(v) \approx v$ or $u \approx v \sigma(v)^{*}$.

We now turn to the second property of the one-cocycle property. In the above proof we should note that the unitary $v$ essentially resides on $Z_{1} \cup Z_{2} \cup Z_{3}$ (or belong to $B\left(Z_{1} \cup Z_{2} \cup Z_{3}\right)$ ). We have a total control of $Z_{2}$ (where $\left(e_{2, i}^{(k)}\right)$ resides) and $Z_{3}$ (where $\left(e_{3, i}^{(k)}\right)$ resides); i.e. by applying $\sigma$ many times we can make $Z_{2}, Z_{3}$ as far away as we want. Hence, we only have to control $Z_{1}$, where $f$ resides.

This amounts to showing that for any finite subset $F$ of $\mathbb{Z}$, there is a finite subset $G$ of $\mathbb{Z}$ such that for any $u \in B(\mathbb{Z} \backslash G)$, there are unitaries $u^{\prime} \in Q\left(J_{M}\right)+1$ and $v \in B\left(Z_{1}^{\prime}\right)$ such that $\left\|v^{*} u \sigma(v)-u^{\prime}\right\|<\epsilon$, where $Z_{1}^{\prime}$ must be a finite subset of $\mathbb{Z}$ disjoint from $F$. If we assume that $u \in B\left(Z_{1}^{\prime \prime}\right)$ with a finite subset $Z_{1}^{\prime \prime}$ disjoint from $G$, then $u^{\prime}$ essentially belong to $B\left(Z_{1}\right)$, where $Z_{1}=Z_{1}^{\prime} \cup Z_{1}^{\prime \prime}$ is disjoint from $F$. (We then apply the above arguments to $u^{\prime}$ instead of $u$.)

To get such $u^{\prime}$ we apply the second condition of the one-cocycle property for the ( $n-1$ )-shift a finite number of times, as discussed in the beginning of this proof. Hence, it is indeed possible.

We now turn to the third property of the one-cocycle property. In the above proof we have chosen four partial isometries $b, b_{1}, b_{2}, b_{3}$, which we now have to choose more carefully, i.e. to make them almost commute with a prescribed finite subset of $B$ requiring some commutativity condition on the unitary $u$ and the projection $f$, and the cycles $e_{j, i}^{(k)}$.

To find a partial isometry $b$ onto $Q\left(1-E_{2, k}\right)$ from a subprojection of $Q\left(e_{2,0}^{(k)} e_{3,0}^{(k)}\right)$ we have derived the condition

$$
\left[Q\left(f\left(1-E_{2, k}\right)\right)\right] \leq\left[Q\left(e_{2,0}^{(k)} e_{3,0}^{(k)}\right)\right]
$$

Since we have to impose the condition that $b$ should commute with a prescribed finitedimensional $C^{*}$-subalgebra $B_{1}$ of $B$, we have to assume that $u, f, e_{2, i}^{(k)}$, and $e_{3, i}^{(k)}$ almost commute with $B_{1}$ and moreover that

$$
\left[Q\left(f\left(1-E_{2, k}\right)\right) P\right] \leq\left[Q\left(e_{2,0}^{(k)} e_{3,0}^{(k)}\right) P\right]
$$

for all minimal central projections $P$ in $B_{1}$.
The first part of the condition is easy. The near commutativity of $e_{2, i}^{(k)}$ and $e_{3, i}^{(k)}$ with $B_{1}$ can be assumed just by shifting them by using $\sigma$. The near commutativity of $f$ can also be assumed by making $f$ large (in the ideal retaining the condition $\sigma(f) \approx f$ ). The near commutativity of $u$ follows from the induction assumption, where we use the stable onecocycle property for the $(n-1)$-shift a finite number of times. Hence, we concentrate on the second part.

By lifting $B_{1}$ to a finite-dimensional $C^{*}$-subalgebra of $A_{n+m}$, we regard $B_{1}$ as a $C^{*}$-subalgebra of $A_{n+m}\left(Z_{0}\right)$ for some finite subset $Z_{0}$ of $\mathbb{Z}$. We can assume that $u, f \in B_{1}^{\prime}$ and $Z_{2} \cup Z_{3}$ is disjoint from $Z_{0}$. From the latter it follows that

$$
\left[e_{2,0}^{(k)} e_{3,0}^{(k)} p\right]=\left[e_{2,0}^{(k)} e_{3,0}^{(k)}\right][p] \geq C^{2}\left(x_{0} \cdots x_{n-1}\right)^{M_{2}+M_{3}}[p]
$$

for a minimal central projection $p$ in $B_{1}$.
Let $M_{0}=\left|Z_{0}\right|$ and let $p$ be a minimal central projection in $B_{1}$. Let $p=\sum_{v \in V^{\prime}} p_{v}$, where $p_{v}$ is a non-zero projection such that $\left[p_{v}\right]=c_{v} x^{v}$ with $c_{v} \in \mathbb{N}$ and $V^{\prime}$ is a subset of $V_{M_{0}}$. Suppose that $f \in J_{M}$ for some $M \in \mathbb{N}$ with $M>M_{0}$. Since $f p_{v}$ belongs to the ideal $J_{M} \cap I(v)=I\left(\sum_{i=0}^{n-1} M e_{i}+\sum_{i=n}^{n+m-1} v_{i} e_{i}\right)$, one can estimate that

$$
\left[f\left(1-E_{2, k}\right) p\right](x) \leq[f p](x) \leq\left(x_{0} x_{1} \cdots x_{n-1}\right)^{M} \sum_{v \in V^{\prime}} D_{v} x^{v}[n]
$$

on $\Lambda_{n+m}$ for some $D_{v} \geq 0$, where $x^{v}[n]=x_{n}^{v_{n}} \cdots x_{n+m-1}^{v_{n+m-1}}$.
Set $M=M_{0}+M_{2}+M_{3}+1$. Since

$$
\left[e_{2,0}^{(k)} e_{3,0}^{(k)} p\right](x) \geq C^{2}\left(x_{0} \cdots x_{n-1}\right)^{M_{0}+M_{2}+M_{3}} \sum_{v \in V^{\prime}} c_{v} x^{v}[n]
$$

on $\Lambda_{n+m}$, it follows that

$$
\left\{x \in \Lambda_{n+m} \mid\left[f\left(1-E_{2, k}\right) p\right](x)<\left[e_{2,0}^{(k)} e_{3,0}^{(k)} p\right](x)\right\}
$$

includes a neighborhood of the vertex $x_{i}=1$, excluding the vertex itself, for $i=$ $0, \ldots, n-1$. Hence, for a sufficiently large $k$ we have that for all minimal central projections $p$ in $B_{1}$,

$$
\left[f\left(1-E_{2, k}\right) p\right](x)<\left[e_{2,0}^{(k)} e_{3,0}^{(k)} p\right](x)
$$

on $\Lambda_{n+m}^{(1 / 2)}$.
As in the proof of Lemma 4.11 we choose $C_{n+i} \in \mathbb{N}$ such that

$$
\sum_{i=0}^{m-1} C_{n+i} x_{n+i}^{k_{n+i}}-\left(x_{0} \cdots x_{n-1}\right)^{-M}\left[f\left(1-E_{2, k}\right) p\right]
$$

is strictly positive on $\Lambda_{n+m} \backslash \Lambda_{n+m}^{(1 / 2)}$. Hence,

$$
\left[e_{2,0}^{(k)} e_{3,0}^{(k)} p\right]+q-\left[f\left(1-E_{2, k}\right) p\right]
$$

is strictly positive on $\Lambda_{n+m}$ except for the vertices $x_{0}=1, \ldots, x_{n-1}=1$, where $q=\left(x_{0} \cdots x_{n-1}\right)^{M} \sum_{i=0}^{m-1} C_{n+i} x_{n+i}^{k_{n+i}}$. We express each term as linear combinations of $x^{v}, v \in V_{L}$ for a sufficiently large $L \in \mathbb{N}$. We have to show that the coefficient is positive on each point of the $i$-vanguard $V_{L}^{i}$ of $\left[e_{2,0}^{(k)} e_{3,0}^{(k)} p\right]+q-\left[f\left(1-E_{2, k}\right) p\right]$ for $i=0, \ldots, n-1$.

Let $v \in V_{L}^{i}$. If $v_{j}<M$ for some $j \neq i$ in $\{0,1, \ldots, n-1\}$, the contribution from [ $f\left(1-E_{2, k}\right) p$ ] must be zero; hence the coefficient is positive.

Suppose that $v_{j} \geq M$ for all $j \neq i$ in $\{0, \ldots, n-1\}$. If $v_{\ell} \geq k_{\ell}$ for some $\ell \geq n$, then $w=\sum_{j \neq i ; j<n} M e_{j}+\left(L-(n-1) M-k_{\ell}\right) e_{i}+k_{\ell} e_{\ell}$ is in front of $v$ and the contribution to the coefficient of $x^{w}$ from $q-\left[f\left(1-E_{2, k}\right) p\right]$, and hence from $\left[e_{2,0}^{(k)} e_{3,0}^{(k)} p\right]+q-\left[f\left(1-E_{2, k}\right) p\right]$,
must be positive, which entails that $w=v$ with positive coefficient of $x^{v}$. So we now have to consider the case $v_{\ell}<k_{\ell}$ for all $\ell \geq n$. In this case the contribution to the coefficient of $x^{v}$ from $q$ is zero. If there is no $w \in V^{\prime}$ such that $D_{w} \neq 0$ and $w_{\ell} \leq v_{\ell}$ for all $\ell \geq n$, then the contribution from $\left[f\left(1-E_{2, k}\right) p\right.$ ] is also zero; thus the coefficient of $x^{v}$ must be positive because it is only contributed from $\left[e_{2,0}^{(k)} e_{3,0}^{(k)} p\right]$. Hence, we are left with the case where there must be $w \in V^{\prime}$ such that $D_{w} \neq 0$ and $w_{\ell} \leq v_{\ell}$ for all $\ell \geq n$. Then there must be $v^{\prime} \in V_{L}$ such that $v_{j}^{\prime}<M$ for $j \neq i$ in $\{0,1, \ldots, n-1\}$ and $v_{\ell}^{\prime}=w_{\ell}$ for all $\ell \geq n$ such that the contribution to the coefficient of $x^{v^{\prime}}$ from $\left[e_{2,0}^{(k)} e_{3,0}^{(k)} p\right]$ is non-zero. Since the contribution to it from $\left[f\left(1-E_{2, k}\right) p\right]$ must be zero, this implies that $v$ cannot be in $V_{L}^{i}$. Thus, we have shown that the coefficients are all positive on the vanguards $V_{L}^{i}$.

Hence, if $C_{n+i}$ are sufficiently large as well as $k$, we can conclude that $\left[e_{2,0}^{(k)} e_{3,0}^{(k)} p\right]+$ $q-\left[f\left(1-E_{2, k}\right) p\right]$ is positive for all minimal central projections $p \in B_{1}$. Thus, it follows that

$$
\left[Q\left(e_{2,0}^{(k)} e_{3,0}^{(k)} p\right)\right]-\left[Q\left(f\left(1-E_{2, k}\right) p\right)\right]
$$

is positive on the quotient $B$ for all minimal central projection $p \in B_{1}$. Hence, we can choose a partial isometry $b$ in the commutant of $B_{1}$ such that

$$
b b^{*}=Q\left(f\left(1-E_{2, k}\right)\right), \quad b^{*} b \leq Q\left(e_{2,0}^{(k)} e_{3,0}^{(k)}\right) .
$$

We then define partial isometries $d, d^{\prime}$ onto $Q\left(f\left(1-E_{2, k}\right)\right)$ and construct $C^{*}$-subalgebras $D, D^{\prime}$ of $B$ (which are isomorphic to $M_{N_{2}+1}$ ) as before.

We now work in the quotient $B=Q\left(A_{n+m}\right)$; so we omit the symbol $Q$. Define

$$
V=\sum_{i=0}^{N_{2}-1}\left(\operatorname{Ad}\left(\zeta^{\prime} x y u\right) \sigma\right)^{i}\left(\left(d^{\prime}\right)^{*}\right)(\operatorname{Ad}(\zeta x y) \sigma)^{i}(d)+f\left(1-E_{2, k}\right) .
$$

Then $V$ is a partial isometry, which we can assume almost commute with a prescribed finite subset of $B$ by imposing such a condition on $u, d, d^{\prime}$. Note that the map $\operatorname{Ad} V$ : $x \mapsto V x V^{*}$ defines an isomorphism of $D$ onto $D_{1}$ and satisfies that $\operatorname{Ad}\left(\zeta^{\prime} x y u\right) \sigma \circ$ $\operatorname{Ad} V \approx \operatorname{Ad} V \circ \operatorname{Ad}(\zeta x y) \sigma$ on $D$, where $\operatorname{Ad}(\zeta x y) \sigma$ (respectively $\operatorname{Ad}\left(\zeta^{\prime} x y u\right) \sigma$ ) leaves $D$ (respectively $D^{\prime}$ ) almost invariant. Thus, if we choose $f_{j, i}, j=1,2$, then we may define $f_{j, i}^{\prime}=V f_{j, i} V^{*}$, which implies that we may set $b_{1}=V f_{1,0}$ and $b_{2}=V f_{2,0}$, which almost commute with a prescribed finite subset.

Since $V^{*} V=1_{D}$, the $N_{2}$-cycle $\left(f_{3, i}\right)$ is defined by

$$
f_{3, i}=e_{2, i}^{(k)}\left(1-V^{*} V\right)
$$

and similarly

$$
f_{3, i}^{\prime}=e_{2, i}^{(k)}\left(1-V V^{*}\right)
$$

Since we can assume that $V$ as well as $e_{2, i}^{(k)}$ almost commutes with elements of a prescribed finite-dimensional $C^{*}$-subalgebra $B_{1}$ of $B$, we have that $\left[f_{3, i} p\right]=\left[f_{3, i}^{\prime} p\right]$ for all minimal central projection $p$ in $B_{1}$. Thus, we can find a partial isometry $b_{3}$ which almost commutes with $B_{1}$ and $b_{3}^{*} b_{3}=f_{3,0}$ and $b_{3} b_{3}^{*}=f_{3,0}^{\prime}$.

Having chosen $b, b_{1}, b_{2}$, and $b_{3}$ as above, the third property also follows as the first property.

Acknowledgement. The author would like to thank one of the referees for calling his attention to $[9,11]$.

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