# A proof of Erdős's $B+B+t$ conjecture 

By Bryna Kra and Joel Moreira and Florian K. Richter and Donald Robertson

January 10, 2024


#### Abstract

We show that every set $A$ of natural numbers with positive upper Banach density can be shifted to contain the restricted sumset $\left\{b_{1}+b_{2}: b_{1}, b_{2} \in B\right.$ and $\left.b_{1} \neq b_{2}\right\}$ for some infinite set $B \subset A$.


## 1. Introduction

One of the most celebrated results in Ramsey theory is van der Waerden's theorem [22]: if one partitions the natural numbers into finitely many pieces, one of those pieces must contain arbitrarily long arithmetic progressions. Often, it transpires that a partition result holds because one of the cells of the partition is large in a suitable sense. In the case of van der Waerden's theorem, positive density was conjectured by Erdős and Turán [5] to guarantee the existence of arbitrarily long progressions. After work of Roth [20] established the case of length three progressions, Szemerédi settled the conjecture positively in general [21], showing that any set $A \subset \mathbb{N}$ with positive upper density contains arbitrarily long arithmetic progressions.

Contemporaneously, Hindman [12] proved a landmark result involving infinite arithmetic patterns: for every finite partition of the natural numbers there is an infinite set $I \subset \mathbb{N}$ such that

$$
\left\{\sum_{i \in F} i: F \subset I, 0<|F|<\infty\right\}
$$

is in one cell of the partition. Looking to connect the two major achievements - Hindman's theorem and Szemerédi's theorem - Erdős formulated the following conjecture on multiple occasions.

Conjecture 1.1 (Erdős [6, Page 305], [7, Pages 57-58], and [8, Page 105]). For any $A \subset \mathbb{N}$ with positive density there exists an infinite set $B \subset A$ and a number $t \in \mathbb{N}$ such that

$$
A-t \supset\left\{b_{1}+b_{2}: b_{1}, b_{2} \in B \text { and } b_{1} \neq b_{2}\right\} .
$$

This problem was studied by various authors, including Nathanson [19], Kazhdan (see [19, 13]), and Hindman [13, Section 11]. Hindman provided several equivalent forms, including a natural reformulation using the Stone-Čech compactification of the integers. A special case of Conjecture 1.1, also conjectured by Erdős, was resolved in [18], asserting that, under the same assumptions, $A$ contains a sumset

$$
B+C=\{b+c: b \in B, c \in C\}
$$

of two infinite sets $B, C \subset \mathbb{N}$. Further recent progress in this direction has been made in $[3,14,16]$, and further history on Conjecture 1.1 and surrounding problems can be found in $[13,17,18,19]$.

Our main theorem resolves Conjecture 1.1. To state our result precisely, recall that a Følner sequence $\Phi$ on $\mathbb{N}$ is any sequence $N \mapsto \Phi_{N}$ of finite subsets of $\mathbb{N}$ with the property that

$$
\lim _{N \rightarrow \infty} \frac{\left|\Phi_{N} \cap\left(\Phi_{N}+t\right)\right|}{\left|\Phi_{N}\right|}=1
$$

for all $t \in \mathbb{N}$. A set $A \subset \mathbb{N}$ has positive upper Banach density if

$$
\lim _{N \rightarrow \infty} \frac{\left|A \cap \Phi_{N}\right|}{\left|\Phi_{N}\right|}>0
$$

for some Følner sequence $\Phi$.
Theorem 1.2. For any $A \subset \mathbb{N}$ with positive upper Banach density, the following hold:
(i) There exist an infinite set $B \subset A$ and a shift $t \in \mathbb{N}$ such that

$$
\left\{b_{1}+b_{2}: b_{1}, b_{2} \in B \text { and } b_{1} \neq b_{2}\right\} \subset A-t
$$

(ii) There exist an infinite set $B \subset \mathbb{N}$ and a shift $t \in \mathbb{N}$ such that

$$
B \cup\left\{b_{1}+b_{2}: b_{1}, b_{2} \in B \text { and } b_{1} \neq b_{2}\right\} \subset A-t
$$

Note that in the formulation of Theorem 1.2 it is not possible to omit the shift by $t$ or remove the condition $b_{1} \neq b_{2}$ in either conclusion (see the discussion in [18] after Question 6.2). Also, it was observed by Hindman in [13] that writing $t=2 r+s$ for $r \in \mathbb{N}$ and $s \in\{0,1\}$ and replacing $B$ by $B-r$, one obtains the following corollary from Theorem 1.2.

Corollary 1.3. For any set $A$ of even integers with positive upper Banach density there exists an infinite set $B \subset \mathbb{N}$ such that $A \supset\left\{b_{1}+b_{2}: b_{1}, b_{2} \in B\right.$ and $\left.b_{1} \neq b_{2}\right\}$.

Our proof of Theorem 1.2 uses ergodic theory and builds on the new dynamical methods developed in [16] to find infinite patterns in sets with positive upper density. To formulate our main dynamical result we recall some basic terminology. By a topological system, we mean a pair $(X, T)$ where $X$ is a compact metric space and $T: X \rightarrow X$ is a homeomorphism. A system is a triple $(X, \mu, T)$, where $(X, T)$ is a topological system and $\mu$ is a $T$-invariant Borel probability measure on $X$. The system is ergodic if any $T$-invariant Borel subset of $X$ has either measure 0 or measure 1, and equivalently we say
that $\mu$ is ergodic for $T$. Given a system $(X, \mu, T)$, a point $a \in X$ is generic for $\mu$ along a Følner sequence $\Phi$, written $a \in \operatorname{gen}(\mu, \Phi)$, if

$$
\mu=\lim _{N \rightarrow \infty} \frac{1}{\left|\Phi_{N}\right|} \sum_{n \in \Phi_{N}} \delta_{T^{n} a}
$$

where $\delta_{x}$ is the Dirac measure at $x \in X$ and the limit is in the weak* topology. This allows us to formulate a dynamical result equivalent to Theorem 1.2.

Theorem 1.4. Let $(X, \mu, T)$ be an ergodic system, let $a \in \operatorname{gen}(\mu, \Phi)$ for some Følner sequence $\Phi$, and let $E \subset X$ be an open set with $\mu(E)>0$.
(i) There exist $x_{1}, x_{2} \in X, t \in \mathbb{N}$, and a strictly increasing sequence $n_{1}<n_{2}<\ldots$ of integers such that $x_{1} \in E, T^{t} x_{2} \in E$, and $(T \times T)^{n_{i}}\left(a, x_{1}\right) \rightarrow\left(x_{1}, x_{2}\right)$ as $i \rightarrow \infty$.
(ii) There exist $x_{1}, x_{2} \in X, t \in \mathbb{N}$, and a strictly increasing sequence $n_{1}<n_{2}<\ldots$ of integers such that $(T \times T)^{t}\left(x_{1}, x_{2}\right) \in E \times E$ and $(T \times T)^{n_{i}}\left(a, x_{1}\right) \rightarrow\left(x_{1}, x_{2}\right)$ as $i \rightarrow \infty$.

A proof of the equivalence between Theorem 1.2 and Theorem 1.4 is given in Section 2, and the proof of Theorem 1.4 is given in Section 3. For a comparison between the techniques in [16] and this paper, and an outline of how the new difficulties arising are overcome, see Section 3.1.

We conclude the introduction with a natural conjecture on a higher order version of our main theorem.

Conjecture 1.5. Let $A \subset \mathbb{N}$ have positive upper Banach density and let $k \in \mathbb{N}$. Then there exist an infinite set $B \subset \mathbb{N}$ and a shift $t \in \mathbb{N}$ such that

$$
\begin{equation*}
A-t \supset\left\{\sum_{n \in F} n: F \subset B, 0<|F|<k\right\} . \tag{1.1}
\end{equation*}
$$

We remark that an example of Straus answering an earlier question of Erdős (see [2, Theorem 2.2] and [13, Theorem 11.6]) shows that $k$ can not be replaced by infinity in (1.1).

Acknowledgements: BK acknowledges National Science Foundation grant DMS-2054643, JM and FKR thank the organizers of the conference "Ultramath2022" during which part of this project was completed, and DR acknowledges EPSRC grant V050362. We thank the anonymous referee for helpful suggestions and comments.

## 2. Reduction to a dynamical statement

In this section we show the equivalence between Theorems 1.2 and 1.4 , beginning with the easier implication.

Proof that Theorem 1.2 implies Theorem 1.4. We prove that part (i) of Theorem 1.2 im plies part (i) of Theorem 1.4; the same proof with obvious modifications shows that part (ii) of Theorem 1.2 implies part (ii) of Theorem 1.4.

Let $(X, \mu, T)$ be an ergodic system, let $a \in \operatorname{gen}(\mu, \Phi)$ for some Følner sequence $\Phi$, and let $E \subset X$ be an open set with $\mu(E)>0$. Since $E$ is open, there exists some point $y \in E$ that lies in the support of $\mu$. Let $U$ be an open ball centered at $y$ (and so $\mu(U)>0)$ whose closure is contained in $E$. Since $a \in \operatorname{gen}(\mu, \Phi)$, the set $A:=\{n \in$ $\left.\mathbb{N}: T^{n} a \in U\right\}$ has positive upper Banach density. Part (i) of Theorem 1.2 then implies that $A \supset\left\{b_{1}+b_{2}: b_{1}, b_{2} \in B, b_{1} \neq b_{2}\right\}+t$ for some infinite set $B \subset A$ and some $t \in \mathbb{N}$. Compactness of $X$ yields an increasing sequence $\left(n_{i}\right)_{i \in \mathbb{N}}$, taking values in $B$ for which $x_{1}:=\lim _{i \rightarrow \infty} T^{n_{i}} a$ exists. Passing to a subsequence of $\left(n_{i}\right)_{i \in \mathbb{N}}$ if needed, the limit $x_{2}:=\lim _{i \rightarrow \infty} T^{n_{i}} x_{1}$ also exists. Since $B \subset A$, it follows that $x_{1} \in \bar{U} \subset E$. Since $\left\{b_{1}+b_{2}: b_{1}, b_{2} \in B, b_{1} \neq b_{2}\right\}+t \subset A$, we have that $T^{t} x_{2} \in E$ as well.

When the set $A$ in Theorem 1.2 is of the form

$$
A=\left\{n \in \mathbb{N}:\|\theta+n \alpha\|_{\mathbb{R} / \mathbb{Z}}<\varepsilon\right\}
$$

for some $\varepsilon>0$ (or, more generally, is a Bohr set), the existence of a set $B \subset \mathbb{N}$ satisfying the conclusion of Theorem 1.2 is connected to the behavior of 3 -term arithmetic progressions $\theta, \theta+\beta, \theta+2 \beta$ in $\mathbb{R} / \mathbb{Z}$ (or, more generally, in the underlying group). For arbitrary $A \subset \mathbb{N}$ we bridge the gap between the combinatorial statement Theorem 1.2 and the dynamical statement Theorem 1.4 using a dynamical variant of 3 -term progressions defined as follows.

Definition 2.1. Given a topological system $(X, T)$, a point $\left(x_{0}, x_{1}, x_{2}\right) \in X^{3}$ is called a (3-term) Erdős progression if there exists a strictly increasing sequence $n_{1}<n_{2}<\cdots$ of integers such that $(T \times T)^{n_{i}}\left(x_{0}, x_{1}\right) \rightarrow\left(x_{1}, x_{2}\right)$ as $i \rightarrow \infty$.

The role played in this paper by Erdős progressions parallels the role played by Erdős cubes in [15]. Various other notions of dynamical progressions, for example those in [9, $15,11]$, have already been used for related questions, but the one we use does not seem to have been defined previously. We remark that in group rotations all the notions of dynamical progressions agree with the conventional notion of arithmetic progression.

The next result completes the translation between ergodic theory and combinatorics by connecting Erdős progressions and sumsets.

Theorem 2.2. Fix a topological system $(X, T)$ and open sets $U, V \subset X$. If there exists an Erdős progression $\left(x_{0}, x_{1}, x_{2}\right) \in X^{3}$ with $x_{1} \in U$ and $x_{2} \in V$, then there exists some infinite set $B \subset\left\{n \in \mathbb{N}: T^{n} x_{0} \in U\right\}$ such that $\left\{b_{1}+b_{2}: b_{1}, b_{2} \in B, b_{1} \neq b_{2}\right\}$ is a subset of $\left\{n \in \mathbb{N}: T^{n} x_{0} \in V\right\}$.

Proof. Let $c: \mathbb{N} \rightarrow \mathbb{N}$ be a strictly increasing sequence such that $(T \times T)^{c(n)}\left(x_{0}, x_{1}\right) \rightarrow$ $\left(x_{1}, x_{2}\right)$. Since $U$ is a neighborhood of $x_{1}$, by refining the sequence $c(n)$ we can assume without loss of generality that $\{c(n): n \in \mathbb{N}\} \subset\left\{n \in \mathbb{N}: T^{n} x_{0} \in U\right\}$.

We now construct the set $B \subset\{c(n): n \in \mathbb{N}\}$ inductively. First choose $b(1)$ in $\{c(i): i \in \mathbb{N}\}$ with $T^{b(1)} x_{1} \in V$. Note that with this choice of $b(1)$ the set $\left(T^{-b(1)} V\right) \times V$ is a neighborhood of $\left(x_{1}, x_{2}\right)$. Next, choose $b(2)$ in $\{c(i): i \in \mathbb{N}\}$ with $b(2)>b(1)$ and

$$
(T \times T)^{b(2)}\left(x_{0}, x_{1}\right) \in\left(T^{-b(1)} V\right) \times V
$$

It follows that $T^{b(1)+b(2)} x_{0} \in V$ and $x_{1} \in T^{-b(2)} V \cap T^{-b(1)} V$.
Supposing that, by induction, we have found $b(1)<\cdots<b(n) \subset\{c(n): n \in \mathbb{N}\}$ with

$$
x_{0} \in \bigcap_{1 \leq i<j \leq n} T^{-b(i)-b(j)} V \quad \text { and } \quad x_{1} \in \bigcap_{1 \leq i \leq n} T^{-b(i)} V,
$$

we choose $b(n+1) \in\{c(i): i \in \mathbb{N}\}$ with $b(n+1)>b(n)$ and

$$
(T \times T)^{b(n+1)}\left(x_{0}, x_{1}\right) \in\left(\bigcap_{1 \leq i \leq n} T^{-b(i)} V\right) \times V
$$

This is possible because

$$
\left(\bigcap_{1 \leq i \leq n} T^{-b(i)} V\right) \times V
$$

is a neighborhood of $\left(x_{1}, x_{2}\right)$ and $(T \times T)^{c(n)}\left(x_{0}, x_{1}\right) \rightarrow\left(x_{1}, x_{2}\right)$ as $n \rightarrow \infty$. Together with the inductive hypothesis, this implies

$$
x_{0} \in \bigcap_{1 \leq i<j \leq n+1} T^{-b(i)-b(j)} V \quad \text { and } \quad x_{1} \in \bigcap_{1 \leq i \leq n+1} T^{-b(i)} V
$$

concluding the induction. Taking $B=\{b(i): i \in \mathbb{N}\}$ finishes the proof.
To deduce Theorem 1.2 from Theorem 1.4 using Theorem 2.2, we use the following version of the Furstenberg correspondence principle.

Proposition 2.3 ([16, Theorem 2.10]). Given a set $A \subset \mathbb{N}$ with positive upper Banach density there exists an ergodic system $(X, \mu, T)$, a Følner sequence $\Phi$, a point $a \in \operatorname{gen}(\mu, \Phi)$, and a clopen set $E \subset X$ such that $\mu(E)>0$ and $A=\left\{n \in \mathbb{N}: T^{n} a \in E\right\}$.

Proof that Theorem 1.4 implies Theorem 1.2. Suppose $A \subset \mathbb{N}$ has positive upper Banach density. Invoking Proposition 2.3 we find an ergodic system $(X, \mu, T)$, a point $a \in \operatorname{gen}(\mu, \Phi)$, a Følner sequence $\Phi$, and a clopen set $E \subset X$ such that $\mu(E)>0$ and $A=\left\{n \in \mathbb{N}: T^{n} a \in E\right\}$. Using Theorem 1.4, part (i), we can find $t \in \mathbb{N}$ and an Erdős progression of the form $\left(a, x_{1}, x_{2}\right) \in X^{3}$ such that $x_{1} \in E$ and $x_{2} \in T^{-t} E$. It now follows from Theorem 2.2, applied with $U=E$ and $V=T^{-t} E$, that there exists an infinite set $B \subset\left\{n \in \mathbb{N}: T^{n} a \in E\right\}=A$ such that

$$
A-t=\left\{n \in \mathbb{N}: T^{n} a \in T^{-t} E\right\} \supset\left\{b_{1}+b_{2}: b_{1}, b_{2} \in B, b_{1} \neq b_{2}\right\}
$$

completing the proof of Theorem 1.2, part (i). If we invoke part (ii) of Theorem 1.4 instead, then the same argument, but using Theorem 2.2 applied with $U=V=T^{-t} E$, yields Theorem 1.2, part (ii).

## 3. Proof of the dynamical statement (Theorem 1.4)

### 3.1. Outline of the proof

Our first observation is that if $\left(x_{0}, x_{1}\right)$ is generic for a $(T \times T)$-invariant measure on $X \times X$ and $\left(x_{1}, x_{2}\right)$ is in the support of that measure, then $\left(x_{0}, x_{1}, x_{2}\right)$ forms an Erdős progression. One may be tempted, then, to find pairs having these properties with respect to the product measure $\mu \times \mu$ on $X \times X$. However, it may be the case that the product measure is not $T \times T$-ergodic, and so typical pairs may only be generic for one of the ergodic components of $\mu \times \mu$. This possibility leads us to consider the ergodic decomposition of $\mu \times \mu$.

We use the notation $\left(x_{1}, x_{2}\right) \mapsto \lambda_{\left(x_{1}, x_{2}\right)}$ to denote an ergodic decomposition of ( $X \times$ $X, \mu \times \mu, T \times T)$. Since we aim to produce an Erdős progression with prescribed first coordinate $a$, as in [16] we make use of a decomposition $\lambda_{\left(x_{1}, x_{2}\right)}$ that is defined for every pair $\left(x_{1}, x_{2}\right)$ and continuous as a function of $\left(x_{1}, x_{2}\right)$. We show in the next section that, without loss of generality, we may assume that our system admits such a decomposition, referred to as a continuous ergodic decomposition.

To find an Erdős progression $\left(a, x_{1}, x_{2}\right) \in X \times X \times X$, it then suffices to find a pair $\left(x_{1}, x_{2}\right)$ with all of the following properties.

1. The pair $\left(a, x_{1}\right) \in X \times X$ is generic for the measure $\lambda_{\left(a, x_{1}\right)}$.
2. The pair $\left(x_{1}, x_{2}\right) \in X \times X$ lies in the support of $\lambda_{\left(x_{1}, x_{2}\right)}$.
3. The measures $\lambda_{\left(a, x_{1}\right)}$ and $\lambda_{\left(x_{1}, x_{2}\right)}$ are equal.

While Properties 1 and 2 hold for $(\mu \times \mu)$-almost every point $\left(x_{1}, x_{2}\right) \in X \times X$, the explicit construction of a continuous ergodic decomposition in [16] tells us that Property 3 holds if and only if the triple ( $a, x_{1}, x_{2}$ ) sits above a three-term progression in the Kronecker factor of $(X, \mu, T)$. Since $a$ is fixed, this constitutes a set of zero measure with respect to $\mu \times \mu$ whenever the Kronecker factor is not finite.

We must therefore show that Properties 1 and 2 hold within the set of points $\left(x_{1}, x_{2}\right)$ for which ( $a, x_{1}, x_{2}$ ) projects onto a three-term progression in the Kronecker factor. To that end, we introduce a natural measure $\sigma$ on $X \times X$ giving full measure to the set of points $\left(x_{1}, x_{2}\right)$ such that $\left(a, x_{1}, x_{2}\right)$ sits above a three-term progression. Thus $\sigma$ is a measure with the property that almost every pair $\left(x_{1}, x_{2}\right) \in X \times X$ satisfies Property 3. Most of the work then goes into showing that the first two properties hold for $\sigma$-almost every pair $\left(x_{1}, x_{2}\right)$. This is where the present work diverges from [16]. To establish the properties we want, it is necessary to understand in greater detail than [16] the disintegration of $\mu$ over the Kronecker factor.

### 3.2. Using continuous factor maps

Throughout this section, we make use of two types of factor maps from a system $(X, \mu, T)$ to another system $(Y, \nu, S)$.

- Measurable factor maps: a measurable function $\pi: X \rightarrow Y$ such that $\pi(\mu)=\nu$ and $\pi \circ T=S \circ \pi \mu$-almost everywhere.
- Continuous factor maps: a continuous surjection $\pi: X \rightarrow Y$ that such that $\pi(\mu)=\nu$
and $\pi \circ T=S \circ \pi$ everywhere.
If there exists a measurable factor map $\pi: X \rightarrow Y$, then $(Y, \nu, S)$ is called a factor of ( $X, \mu, T$ ).

In his proof of Szemerédi's theorem, Furstenberg [9] shows that in order to understand the behavior of 3-term dynamical progressions, it suffices to consider their projections onto the maximal group rotation factor. We use an analogous method to study 3-term Erdős progressions.

A group rotation is a system of the form $(Z, m, R)$, where $Z$ is a compact abelian group, $m$ is the Haar measure on $Z$, and $R: Z \rightarrow Z$ is a rotation of the form $R(z)=z+\alpha$ for a fixed element $\alpha \in Z$. Whenever $(Z, m, R)$ is a group rotation, we assume that the metric on $Z$ is chosen such that $z \mapsto z+w$ is an isometry for all $w \in Z$.

Every ergodic system $(X, \mu, T)$ possesses a maximal group rotation factor called its Kronecker factor (see [10, Section 3]). In general, the factor map from an ergodic system $(X, \mu, T)$ onto its Kronecker factor $(Z, m, R)$ is only a measurable factor map. The next lemma, however, shows that in many situations one can assume without loss of generality that the factor map onto the Kronecker factor is continuous, and this is key in our proof of Theorem 1.4.

Proposition 3.1 ([16, Proposition 3.20]). Let $(X, \mu, T)$ be an ergodic system and let $a \in \operatorname{gen}(\mu, \Phi)$ for some Følner sequence $\Phi$. Then there exists an ergodic system $(\tilde{X}, \tilde{\mu}, \tilde{T})$, a Følner sequence $\Psi$, a point $\tilde{a} \in \tilde{X}$ and a continuous factor map $\tilde{\pi}: \tilde{X} \rightarrow X$ such that $\tilde{\pi}(\tilde{a})=a$ and $\tilde{a} \in \operatorname{gen}(\tilde{\mu}, \Psi)$ and $(\tilde{X}, \tilde{\mu}, \tilde{T})$ has a continuous factor map to its Kronecker factor.

With the help of Proposition 3.1 we can reduce the proof of Theorem 1.4 to the following special case.

Theorem 3.2. Let $(X, \mu, T)$ be an ergodic system and assume there is a continuous factor map $\pi$ to its Kronecker. Let $a \in \operatorname{gen}(\mu, \Phi)$ for some Følner sequence $\Phi$, and let $E \subset X$ be a Borel set with $\mu(E)>0$.
(i) There exist $t \in \mathbb{N}$ and an Erdős progression of the form $\left(a, x_{1}, x_{2}\right) \in X^{3}$ such that $x_{1} \in E$ and $T^{t} x_{2} \in E$.
(ii) There exist $t \in \mathbb{N}$ and an Erdős progression of the form $\left(a, x_{1}, x_{2}\right) \in X^{3}$ such that $T^{t} x_{1} \in E$ and $T^{t} x_{2} \in E$.

We remark that unlike in Theorem 1.4, in the formulation of Theorem 3.2 we do not require that $E$ is an open set. In fact, this hypothesis is not needed in Theorem 1.4 either, but without assuming openness of $E$, Theorem 1.4 is no longer equivalent to Theorem 1.2.

Proof that Theorem 3.2 implies Theorem 1.4. We only prove that part (i) of Theorem 3.2 implies part (i) of Theorem 1.4. Similar arguments show the implication between part (ii) of Theorem 3.2 and part (ii) of Theorem 1.4.

Let $(X, \mu, T)$ be an ergodic system, let $a \in \operatorname{gen}(\mu, \Phi)$ for some Følner sequence $\Phi$ and let $E \subset X$ be open and have positive measure. Let $(\tilde{X}, \tilde{\mu}, \tilde{T}), \tilde{a}$ and $\tilde{\pi}$ result from an application of Proposition 3.1 and let $\tilde{E}:=\tilde{\pi}^{-1}(E) \subset \tilde{X}$. As $(\tilde{X}, \tilde{\mu}, \tilde{T})$ has a continuous factor map to its Kronecker factor, we can apply Theorem 3.2 to find $t \in \mathbb{N}$ and an

Erdős progression $\left(\tilde{a}, \tilde{x}_{1}, \tilde{x}_{2}\right) \in \tilde{X}^{3}$ with $\tilde{x}_{1} \in \tilde{E}$ and $\tilde{x}_{2} \in T^{-t} \tilde{E}$. It is then immediate that $\left(\tilde{\pi}(\tilde{a}), \tilde{\pi}\left(\tilde{x}_{1}\right), \tilde{\pi}\left(\tilde{x}_{2}\right)\right)$ is an Erdős progression in $X^{3}$ with $\tilde{\pi}\left(\tilde{x}_{1}\right) \in E$ and $\tilde{\pi}\left(\tilde{x}_{2}\right) \in T^{-t}(E)$.

The proof of Theorem 3.2 is deferred to Section 3.5 until after we have developed the necessary tools. We conclude this section by recalling the continuous ergodic decomposition of the product measure $\mu \times \mu$ from [16]. To do so we make use of the following standard disintegration result.

Theorem 3.3 (See [4, Theorem 5.14]). Given a measurable factor map $\pi: X \rightarrow Y$ between systems $(X, \mu, T)$ and $(Y, \nu, S)$, there exists a measurable map $y \mapsto \mu_{y}$ defined on a full measure subset of $Y$ and taking values in the space $\mathcal{M}(X)$ of Borel probability measures on $X$ with the following properties.
(i) For every bounded, measurable function $f: X \rightarrow \mathbb{R}$, the function

$$
y \mapsto \int_{X} f \mathrm{~d} \mu_{y}
$$

is an almost everywhere defined and Borel measurable function on $Y$ satisfying

$$
\int_{D}\left(\int_{X} f \mathrm{~d} \mu_{y}\right) \mathrm{d} \nu(y)=\int_{\pi^{-1}(D)} f \mathrm{~d} \mu
$$

for all Borel sets $D \subseteq Y$.
(ii) For $\nu$-almost every $y \in Y$, we have $\mu_{y}\left(\pi^{-1}(\{y\})\right)=1$.
(iii) Properties (i) and (ii) uniquely determine the map $y \mapsto \mu_{y}$ in the sense that if $y \mapsto \mu_{y}^{\prime}$ is another measurable map from $Y$ to $\mathcal{M}(X)$ with these properties, then $\mu_{y}=\mu_{y}^{\prime}$ for $\nu$-almost every $y \in Y$.
(iv) For almost every $y \in Y$, we have $T \mu_{y}=\mu_{S y}$.

Fix an ergodic system $(X, \mu, T)$. Let $(Z, m, R)$ be its Kronecker factor, and assume that $\pi$ is a continuous factor map from $(X, \mu, T)$ to $(Z, m, R)$. Also fix a disintegration $z \mapsto \eta_{z}$ of $\mu$ with respect to $\pi$. As in [16, Equation (3.10)], for every $\left(x_{1}, x_{2}\right) \in X \times X$ we define

$$
\begin{equation*}
\lambda_{\left(x_{1}, x_{2}\right)}=\int_{Z} \eta_{z+\pi\left(x_{1}\right)} \times \eta_{z+\pi\left(x_{2}\right)} \mathrm{d} m(z) \tag{3.1}
\end{equation*}
$$

on $X \times X$. Note that $\lambda_{\left(x_{1}, x_{2}\right)}$ does not depend on the choice of disintegration $z \mapsto \eta_{z}$. The following properties are proved in [16, Proposition 3.11].

1. The map $\left(x_{1}, x_{2}\right) \mapsto \lambda_{\left(x_{1}, x_{2}\right)}$ is continuous.
2. The map $\left(x_{1}, x_{2}\right) \mapsto \lambda_{\left(x_{1}, x_{2}\right)}$ is a disintegration of $\mu \times \mu$ in the sense that

$$
\int_{X \times X} \lambda_{\left(x_{1}, x_{2}\right)} \mathrm{d}(\mu \times \mu)\left(x_{1}, x_{2}\right)=\mu \times \mu
$$

holds.
3. For $(\mu \times \mu)$-almost every $\left(x_{1}, x_{2}\right)$, the point $\left(x_{1}, x_{2}\right)$ is generic for $\lambda_{\left(x_{1}, x_{2}\right)}$ and $\lambda_{\left(x_{1}, x_{2}\right)}$ is $T \times T$ ergodic.
4. For every $\left(x_{1}, x_{2}\right) \in X \times X$, we have that $\lambda_{\left(x_{1}, x_{2}\right)}=\lambda_{\left(T x_{1}, T x_{2}\right)}$.

### 3.3. The measure on Erdős progressions

In this section, we introduce a measure $\sigma$ on $X \times X$ that helps us study Erdős progressions beginning at a fixed point $a \in X$.

Fix an ergodic system $(X, \mu, T)$ and a point $a \in X$, let $(Z, m, R)$ denote its Kronecker factor, and further assume that there is a continuous factor map $\pi: X \rightarrow Z$. Moreover, we fix a disintegration $z \mapsto \eta_{z}$ of $\mu$ over $\pi$ as guaranteed by Theorem 3.3.
Definition 3.4. We define the measure

$$
\begin{equation*}
\sigma=\int_{Z} \eta_{z} \times \eta_{2 z-\pi(a)} \mathrm{d} m(z)=\int_{Z} \eta_{\pi(a)+z} \times \eta_{\pi(a)+2 z} \mathrm{~d} m(z) \tag{3.2}
\end{equation*}
$$

on $X \times X$.
For the remainder of this paper, we use $\sigma$ to denote the measure defined by (3.2). Note that the second equality in (3.2) follows from translation invariance of $m$. We stress that $\sigma$ does not depend on the exact choice of disintegration $z \mapsto \eta_{z}$ since any two choices agree $m$-almost everywhere.

The motivation for this definition is that $\sigma$ is a relatively independent joining, putting as unbiased as possible a measure on the set of pairs $\left(x_{1}, x_{2}\right) \in X \times X$ such that

$$
\left(\left(\pi(a), \pi\left(x_{1}\right), \pi\left(x_{2}\right)\right)\right.
$$

forms a 3-term arithmetic progression in $Z$. The connection to three-term progressions is made apparent by the equality

$$
\pi\left(x_{2}\right)-\pi\left(x_{1}\right)=\pi\left(x_{1}\right)-\pi(a)
$$

which holds for $\sigma$-almost every $\left(x_{1}, x_{2}\right)$ and guarantees via (3.1) that $\lambda_{\left(a, x_{1}\right)}=\lambda_{\left(x_{1}, x_{2}\right)}$.
We conclude this section with some lemmas that are of use in the next sections.
Lemma 3.5. Let $\pi_{1}: X \times X \rightarrow X$ denote the projection $\left(x_{1}, x_{2}\right) \mapsto x_{1}$ onto the first coordinate. Then $\pi_{1} \sigma=\mu$.

Proof. For any $f \in C(X)$, we have

$$
\begin{aligned}
\int_{X} f \mathrm{~d}\left(\pi_{1} \sigma\right) & =\int_{X \times X}(f \otimes 1) \mathrm{d} \sigma \\
& =\int_{Z}\left(\int_{X \times X}(f \otimes 1) \mathrm{d}\left(\eta_{z} \times \eta_{2 z-\pi(t)}\right)\right) \mathrm{d} m(z) \\
& =\int_{Z}\left(\int_{X} f \mathrm{~d} \eta_{z}\right) \mathrm{d} m(z)=\int_{X} f \mathrm{~d} \mu,
\end{aligned}
$$

as desired.
Lemma 3.6. Let $\pi_{2}: X \times X \rightarrow X$ denote the projection $\left(x_{1}, x_{2}\right) \mapsto x_{2}$ onto the second coordinate. Then $\frac{1}{2}\left(\pi_{2} \sigma+T \pi_{2} \sigma\right)=\mu$.

Proof. Denote by $2 Z$ the subgroup $\{z+z: z \in Z\}$ and let $\xi$ denote its Haar measure. Ergodicity of $R$ ensures that $Z=(2 Z) \cup R(2 Z)$ and that $m=\frac{1}{2}(\xi+R \xi)$. In particular,
for each $s \in X$ there exists $w \in Z$ such that either $\pi(s)=2 w$ or $\pi(s)=R(2 w)$. In the first case

$$
\pi_{2} \sigma=\int_{Z} \eta_{2(w+z)} \mathrm{d} m(z)=\int_{Z} \eta_{2 z} \mathrm{~d} m(z)=\int_{2 Z} \eta_{u} \mathrm{~d} \xi(u)
$$

and in the second

$$
\pi_{2} \sigma=\int_{Z} \eta_{2(w+z)+\alpha} \mathrm{d} m(z)=\int_{Z} \eta_{2 z+\alpha} \mathrm{d} m(z)=\int_{2 Z+\alpha} \eta_{u} \mathrm{~d}(R \xi)(u)
$$

Since $T \eta_{u}=\eta_{R u}$ and $R^{2} \xi=\xi$, it follows that in either case

$$
\frac{1}{2}\left(\pi_{2} \sigma+T \pi_{2} \sigma\right)=\int_{Z} \eta_{z} \mathrm{~d} \frac{1}{2}(\xi+R \xi)(z)=\mu .
$$

### 3.4. The support of the measure

We maintain the notation of Section 3.3 and assume that $(X, \mu, T)$ is an ergodic system with Kronecker factor $(Z, m, R)$, continuous factor map $\pi:(X, \mu, T) \rightarrow(Z, m, R)$, the measures $\lambda_{x_{1}, x_{2}}$ are those defined in (3.1), and $\sigma$ denotes the measure defined in (3.4). We continue to use the fixed disintegration $z \mapsto \eta_{z}$ of $\mu$ with respect to $\pi$.

Lemma 3.7. For $\sigma$-almost every $\left(x_{1}, x_{2}\right) \in X \times X$, the measures $\lambda_{\left(a, x_{1}\right)}$ and $\lambda_{\left(x_{1}, x_{2}\right)}$ are equal.

Proof. Consider the set

$$
P:=\left\{\left(x_{1}, x_{2}\right) \in X \times X: \pi\left(x_{1}\right)=\pi(a)+z, \pi\left(x_{2}\right)=\pi(a)+2 z \text { for some } z \in Z\right\} .
$$

Combining (3.2) and property (ii) of Theorem 3.3 for the disintegration $z \mapsto \eta_{z}$, it follows that $\sigma(P)=1$ and each $\left(x_{1}, x_{2}\right) \in P$ satisfies

$$
\pi\left(x_{2}\right)-\pi\left(x_{1}\right)=\pi\left(x_{1}\right)-\pi(a)
$$

Thus we have $\lambda_{\left(x_{1}, x_{2}\right)}=\lambda_{\left(a, x_{1}\right)}$ by the defining formula (3.1) and translation invariance of $m$.

Let $\operatorname{supp}(\nu)$ denote the support of a Borel measure $\nu$ and let $\mathcal{F}(X)$ denote the family of closed, nonempty subsets of a given compact metric space $(X, d)$. We endow $\mathcal{F}(X)$ with the Haudsorff metric H, defined by

$$
\mathrm{H}(F, G)=\max \left\{\sup _{x \in F} d(x, G), \sup _{y \in G} d(y, F)\right\}
$$

whenever $F, G \in \mathcal{F}(X)$.
Lemma 3.8. Let $W$ be a compact metric space, $\mathcal{M}(W)$ the space of Borel probability measures on $W$ endowed with the weak* topology, and $\mathcal{F}(W)$ the space of closed, nonempty subsets of $W$ endowed with the Hausdorff metric.

1. The map $\nu \mapsto \operatorname{supp}(\nu)$ from $\mathcal{M}(W)$ to $\mathcal{F}(W)$ is Borel measurable.
2. If $x \mapsto \rho_{x}$ is a measurable map from $W$ to $\mathcal{M}(W)$, then $\left\{x \in W: x \in \operatorname{supp}\left(\rho_{x}\right)\right\}$ is a Borel set.

## Proof.

1. Combining Theorem 17.14, Lemma 17.5, and Theorem 18.9 in [1], the result follows.
2. The map $\psi_{1}(x)=\{x\}$ from $W$ to $\mathcal{F}(W)$ is continuous and hence measurable. By part 1, the map $\psi_{2}(x)=\operatorname{supp}\left(\rho_{x}\right)$ from $W$ to $\mathcal{F}(W)$ is also measurable. Thus $\psi(x)=\left(\psi_{1}(x), \psi_{2}(x)\right)$ from $W$ to $\mathcal{F}(W) \times \mathcal{F}(W)$ is measurable. The set $\Omega=$ $\left\{\left(F_{1}, F_{2}\right) \in \mathcal{F}(W) \times \mathcal{F}(W): F_{1} \cap F_{2} \neq \varnothing\right\}$ is closed, and therefore $\{x \in W: x \in$ $\left.\operatorname{supp}\left(\rho_{x}\right)\right\}=\psi^{-1}(\Omega)$ is Borel.

Lemma 3.9. The disintegration $z \mapsto \eta_{z}$ satisfies $\mu\left(\left\{x \in X: x \in \operatorname{supp}\left(\eta_{\pi(x)}\right)\right\}\right)=1$.
Proof. Write $G=\left\{x \in X: x \in \operatorname{supp}\left(\eta_{\pi(x)}\right)\right\}$, which is Borel measurable by Lemma 3.8, part 2. Since

$$
\mu(G)=\int_{Z} \eta_{z}(G) \mathrm{d} m(z)
$$

it suffices to show $\eta_{z}(G)=1$ for $m$-almost every $z \in Z$. By Theorem 3.3, part (ii), for almost every $z \in Z$ we have $\eta_{z}\left(\pi^{-1}(z)\right)=1$. If $\eta_{z}\left(\pi^{-1}(z)\right)=1$, then $\operatorname{supp}\left(\eta_{z}\right) \subset \pi^{-1}(z)$ because continuity of $\pi$ gives that $\pi^{-1}(z)$ is a closed set, and therefore it is a closed set of full measure. Thus, for $m$-almost every $z \in Z$, we have $\operatorname{supp}\left(\eta_{z}\right) \subset \pi^{-1}(z)$ and hence $\operatorname{supp}\left(\eta_{z}\right) \subset G$. Since $\operatorname{supp}\left(\eta_{z}\right) \subset G$, we have $\eta_{z}(G) \geq \eta_{z}\left(\operatorname{supp}\left(\eta_{z}\right)\right)=1$ for $m$-almost every $z \in Z$.

Write

$$
\begin{equation*}
S=\left\{\left(x_{1}, x_{2}\right) \in X \times X:\left(x_{1}, x_{2}\right) \in \operatorname{supp}\left(\lambda_{\left(x_{1}, x_{2}\right)}\right)\right\} \tag{3.3}
\end{equation*}
$$

and note that part 2 in Lemma 3.8, together with continuity of $\left(x_{1}, x_{2}\right) \mapsto \lambda_{\left(x_{1}, x_{2}\right)}$, implies that $S$ is a Borel subset of $X \times X$. Our goal for the remainder of this section is to show that $\sigma(S)=1$ for every $s \in X$ (see Proposition 3.11).

Proposition 3.10. Fix a system $(X, \mu, T)$ and a continuous factor map $\pi$ to its Kronecker factor $(Z, m, T)$. Also fix a disintegration $z \mapsto \eta_{z}$ over its Kronecker factor ( $Z, m, R$ ). There is a sequence $\delta(j) \rightarrow 0$ such that for almost every $x \in X$ the following holds: for every neighbourhood $U$ of $x$ we have

$$
\begin{equation*}
\lim _{j \rightarrow \infty} \frac{m\left(\left\{z \in Z: \eta_{z}(U)>0\right\} \cap B(\pi(x), \delta(j))\right)}{m(B(\pi(x), \delta(j)))}=1 \tag{3.4}
\end{equation*}
$$

Proof. Consider the map $\Phi: Z \rightarrow \mathcal{F}(X)$ given by $\Phi(z)=\operatorname{supp}\left(\eta_{z}\right)$. This map is Borel measurable by Lemma 3.8 as it is the composition of two Borel measurable functions $z \mapsto \eta_{z}$ and $\nu \mapsto \operatorname{supp}(\nu)$. Applying Lusin's theorem [1, Theorem 12.8] for every $j \in \mathbb{N}$,
there is a closed set $Z_{j} \subset Z$ with $m\left(Z_{j}\right)>1-2^{-j}$ such that $\left.\Phi\right|_{Z_{j}}$ is continuous. By uniform continuity of $\left.\Phi\right|_{Z_{j}}$, there exists a positive number $\delta(j)$ such that for all $z_{1}, z_{2} \in Z_{j}$,

$$
d\left(z_{1}, z_{2}\right) \leqslant \delta(j) \Longrightarrow \mathrm{H}\left(\Phi\left(z_{1}\right), \Phi\left(z_{2}\right)\right)<\frac{1}{j} .
$$

Consider the set

$$
K_{j}=\left\{z \in Z_{j}: m\left(B(z, \delta(j)) \cap Z_{j}\right)>\left(1-\frac{1}{j}\right) m(B(z, \delta(j)))\right\}
$$

for each $j \in \mathbb{N}$. Define

$$
\chi_{j}(z)=\frac{1}{m(B(0, \delta(j)))} \int_{Z_{j}} \mathbf{1}_{B(0, \delta(j))}(w-z) \mathrm{d} m(w)
$$

and note that $\chi_{j}(z) \leqslant 1$ for all $z \in Z$. Since translations on $Z$ are isometries, we have

$$
\begin{equation*}
K_{j}=Z_{j} \cap\left\{z \in Z: \chi_{j}(z)>\left(1-\frac{1}{j}\right)\right\} . \tag{3.5}
\end{equation*}
$$

Using Fubini's theorem, we deduce that

$$
\begin{equation*}
\int_{Z} \chi_{j}(z) \mathrm{d} m(z)=m\left(Z_{j}\right)>1-\frac{1}{2^{j}} \tag{3.6}
\end{equation*}
$$

which combined with $\chi_{j}(z) \leqslant 1$ implies that

$$
\begin{equation*}
m\left(\left\{z \in Z: \chi_{j}(z)>\left(1-\frac{1}{j}\right)\right\}\right) \geqslant 1-\frac{j}{2^{j}} \tag{3.7}
\end{equation*}
$$

Combining (3.5) with (3.6) and (3.7), it follows that $\sum_{j \in \mathbb{N}} m\left(Z \backslash K_{j}\right)<\infty$.
Let

$$
K=\bigcup_{M \geqslant 1} \bigcap_{j \geqslant M} K_{j} .
$$

Observe that, by the Borel-Cantelli lemma, $m(K)=1$. In view of Lemma 3.9, this implies that the set $L:=\left\{x \in X: x \in \operatorname{supp}\left(\eta_{\pi(x)}\right)\right\} \cap \pi^{-1}(K)$ has $\mu(L)=1$. To finish the proof it thus suffices to show that any $x \in L$ satisfies (3.4).

Fix a point $x \in L$ and let $U$ be a neighborhood of $x$. Let $z=\pi(x)$. Since $z \in K$ and $U$ is open, there exists $j_{0} \in \mathbb{N}$ such that for all $j \geqslant j_{0}$ we have $z \in K_{j}$ and $B(x, 1 / j) \subset U$. We claim that for all $j \geqslant j_{0}$, we have

$$
\begin{equation*}
\left.B(z, \delta(j)) \cap Z_{j} \subset H:=\left\{z \in Z: \eta_{z}(U)>0\right)\right\} \tag{3.8}
\end{equation*}
$$

To verify this claim, let $z^{\prime} \in B(z, \delta(j)) \cap Z_{j}$ be arbitrary. Since $\mathrm{H}\left(\Phi(z), \Phi\left(z^{\prime}\right)\right)<1 / j$ and $x \in \Phi(z)$, there exists $x^{\prime} \in \Phi\left(z^{\prime}\right)$ with $d\left(x, x^{\prime}\right)<1 / j$. From $d\left(x, x^{\prime}\right)<1 / j$ it follows that $x^{\prime} \in U$ and using $x^{\prime} \in \Phi\left(z^{\prime}\right)$ we conclude $U \cap \Phi\left(z^{\prime}\right) \neq \varnothing$. Since $\Phi\left(z^{\prime}\right)=\operatorname{supp}\left(\eta_{z^{\prime}}\right)$, it follows that $\eta_{z^{\prime}}(U)>0$ and hence that $z^{\prime} \in H$, proving that (3.8) holds, as claimed.

Since $z \in K_{j}$, it follows from (3.8) and the construction of $K_{j}$ that

$$
\frac{m(H \cap B(z, \delta(j)))}{m(B(z, \delta(j)))} \geqslant 1-\frac{1}{j}
$$

for all $j \geqslant j_{0}$. We conclude that

$$
\lim _{j \rightarrow \infty} \frac{m(H \cap B(z, \delta(j)))}{m(B(z, \delta(j)))}=1
$$

and the proof is complete.
Proposition 3.11. The set $S$ defined in (3.3) satisfies $\sigma(S)=1$.
Proof. Apply Proposition 3.10 to get a sequence $\delta(j) \rightarrow 0$ with the properties therein. Let $L$ denote the set of points satisfying (3.4) which has full $\mu$-measure. We conclude from Lemma 3.5 that $\sigma(L \times X)=1$ and conclude from Lemma 3.6 that

$$
1=\mu(L)=\frac{\sigma(X \times L)+\sigma\left(X \times T^{-1} L\right)}{2}
$$

whence $\sigma(X \times L)=1$. Thus

$$
\sigma(L \times L)=\sigma((X \times L) \cap(L \times X))=1
$$

To prove $\sigma(S)=1$, it therefore suffices to show $L \times L \subset S$.
Let $\left(x_{1}, x_{2}\right) \in L \times L$. Let $U_{1}$ be a neighborhood of $x_{1}$ and let $U_{2}$ be a neighborhood of $x_{2}$. To show $\left(x_{1}, x_{2}\right) \in S$, we have to verify $\lambda_{\left(x_{1}, x_{2}\right)}\left(U_{1} \times U_{2}\right)>0$. For convenience, write $\beta=\pi\left(x_{2}\right)-\pi\left(x_{1}\right)$. By definition,

$$
\lambda_{\left(x_{1}, x_{2}\right)}=\int_{Z} \eta_{z} \times \eta_{z+\beta} \mathrm{d} m(z)
$$

Since $x_{1}, x_{2} \in L$, there exists some $\delta>0$ such that

$$
\begin{equation*}
\frac{\left.m\left(\left\{z \in Z: \eta_{z}\left(U_{1}\right)>0\right)\right\} \cap B\left(\pi\left(x_{1}\right), \delta\right)\right)}{m\left(B\left(\pi\left(x_{1}\right), \delta\right)\right)} \geqslant \frac{3}{4} \tag{3.9}
\end{equation*}
$$

as well as

$$
\begin{equation*}
\frac{\left.m\left(\left\{z \in Z: \eta_{z}\left(U_{2}\right)>0\right)\right\} \cap B\left(\pi\left(x_{2}\right), \delta\right)\right)}{m\left(B\left(\pi\left(x_{2}\right), \delta\right)\right)} \geqslant \frac{3}{4} \tag{3.10}
\end{equation*}
$$

Observe that $\left.\left.\left\{z \in Z: \eta_{z}\left(U_{2}\right)>0\right)\right\}-\beta=\left\{z \in Z: \eta_{z+\beta}\left(U_{2}\right)>0\right)\right\}$, and so (3.10) implies

$$
\begin{equation*}
\frac{\left.m\left(\left\{z \in Z: \eta_{z+\beta}\left(U_{2}\right)>0\right)\right\} \cap B\left(\pi\left(x_{1}\right), \delta\right)\right)}{m\left(B\left(\pi\left(x_{1}\right), \delta\right)\right)} \geqslant \frac{3}{4} \tag{3.11}
\end{equation*}
$$

Define $W=\left\{z \in Z: \eta_{z}\left(U_{1}\right)>0\right.$ and $\left.\left.\eta_{z+\beta}\left(U_{2}\right)>0\right)\right\}$. By (3.9) and (3.11) it follows that $W$ contains at least one-quarter of the ball $B\left(\pi\left(x_{1}\right), \delta\right)$, which implies $m(W)>0$. Since for all $z \in W$ one has

$$
\left(\eta_{z} \times \eta_{z+\beta}\right)\left(U_{1} \times U_{2}\right)>0
$$

and $m(W)>0$, it follows that $\lambda_{\left(x_{1}, x_{2}\right)}\left(U_{1} \times U_{2}\right)>0$ as desired.

### 3.5. Proof of Theorem 3.2

To prove Theorem 3.2 we need one further lemma.
Lemma 3.12 ([16, Lemma 3.18]). Let $(X, \mu, T)$ be an ergodic system, let $a \in \operatorname{gen}(\mu, \Phi)$ for some Følner sequence $\Phi$. Then there exists a Følner sequence $\Psi$ such that for $\mu$-almost every $x_{1} \in X$ the point ( $a, x_{1}$ ) belongs to gen $\left(\lambda_{\left(a, x_{1}\right)}, \Psi\right)$.

Proof. From property (3) after the definition of $\lambda_{(x, y)}$ in (3.1) and Fubini's theorem, there exists (a full measure set of) $b \in \operatorname{supp}(\mu)$ such that for $\mu$-almost every $x \in X$, the point $(b, x)$ is generic for $\lambda_{(b, x)}$ with respect to the Følner sequence $(\{1, \ldots, N\})_{N \in \mathbb{N}}$. Let $\left(G_{j}\right)_{j=1}^{\infty}$ enumerate a countable dense subset of $C(X \times X)$ and, for each $j \in \mathbb{N}$, let $\tilde{G}_{j}(x, y)=\int_{X^{2}} G_{j} \mathrm{~d} \lambda_{(x, y)}$. Since the map $\left(x_{1}, x_{2}\right) \mapsto \lambda_{\left(x_{1}, x_{2}\right)}$ is continuous and $(T \times T)$ invariant, each of the functions $\tilde{G}_{j}$ is also continuous and $(T \times T)$-invariant.

Since $a \in \operatorname{gen}(\mu, \Phi)$ and $b \in \operatorname{supp}(\mu)$, for every $m \in \mathbb{N}$, there exists $s(m) \in \mathbb{N}$ such that $\left\|G_{j}(b, \cdot)-G_{j}\left(T^{s(m)} a, \cdot\right)\right\|_{\infty}<2^{-m}$ for every $j \leqslant m$ and $\left\|\tilde{G}_{j}(b, \cdot)-\tilde{G}_{j}\left(T^{s(m)} a, \cdot\right)\right\|_{\infty}<2^{-m}$ for every $j \leqslant m$. Since $(b, x)$ is generic for $\lambda_{(b, x)}$ for $\mu$-almost every $x \in X$, for each $m \in \mathbb{N}$, there exists some $N(m) \in \mathbb{N}$ for which the continuous function

$$
F_{m}(x):=\max _{1 \leqslant j \leqslant m}\left|\frac{1}{N(m)} \sum_{n=1}^{N(m)} G_{j}\left(T^{n} b, T^{n} x\right)-\tilde{G}_{j}(b, x)\right|
$$

satisfies $\left\|F_{m}\right\|_{L^{1}(\mu)}<1 / 2^{m}$. The choice of $s(m)$ implies that

$$
\tilde{F}_{m}(x):=\max _{1 \leqslant j \leqslant m}\left|\frac{1}{N(m)} \sum_{n=1}^{N(m)} G_{j}\left(T^{n+s(m)} a, T^{n} x\right)-\tilde{G}_{j}\left(T^{s(m)} a, x\right)\right|
$$

satisfies $\left\|\tilde{F}_{m}\right\|_{L^{1}(\mu)}<3 / 2^{m}$. Letting $\Psi_{m}=\{s(m)+1, \ldots, s(m)+N(m)\}$ and using $T \times T$ invariance of $\tilde{G}_{j}$ we deduce that

$$
F_{m}^{\prime}(x):=\tilde{F}_{m}\left(T^{s(m)} x\right)=\max _{1 \leqslant j \leqslant m}\left|\frac{1}{\left|\Psi_{m}\right|} \sum_{n \in \Psi_{m}} G_{j}\left(T^{n} a, T^{n} x\right)-\tilde{G}_{j}(a, x)\right|
$$

Since $\mu$ is $T$-invariant, $\left\|F_{m}^{\prime}\right\|_{L^{1}(\mu)}=\left\|\tilde{F}_{m}\right\|_{L^{1}(\mu)}<3 / 2^{m}$ for every $m \in \mathbb{N}$, hence the function $F(x):=\sum_{m \in \mathbb{N}} F_{m}^{\prime}(x)$ has norm $\|F\|_{L^{1}(\mu)}=\sum\left\|F_{m}^{\prime}\right\|_{L^{1}(\mu)}<\infty$ and is therefore finite almost everywhere. For every point $x_{1} \in X$ for which $F\left(x_{1}\right)<\infty$, we have that $F_{m}^{\prime}\left(x_{1}\right) \rightarrow 0$ as $m \rightarrow \infty$ and hence $\left(a, x_{1}\right) \in \operatorname{gen}\left(\lambda_{\left(a, x_{1}\right)}, \Psi\right)$.

Proof of Theorem 3.2. Fix a system $(X, \mu, T), a \in \operatorname{gen}(\mu, \Phi)$ for some Følner sequence $\Phi$, and $E \subset X$ open with $\mu(E)>0$. Assume $(X, \mu, T)$ has a continuous factor map $\pi$ to its Kronecker factor $(Z, m, R)$. Applying Lemma 3.12 it follows that for $\mu$-almost every $x_{1} \in X$, we have $\left(a, x_{1}\right) \in \operatorname{gen}\left(\lambda_{\left(a, x_{1}\right)}, \Psi\right)$ for some Følner sequence $\Psi$. Since, in view of Lemma 3.5, the projection of $\sigma$ onto the first coordinate equals $\mu$, it follows that for $\sigma$-almost every $\left(x_{1}, x_{2}\right) \in X \times X$ we have $\left(a, x_{1}\right) \in \operatorname{gen}\left(\lambda_{\left(a, x_{1}\right)}, \Psi\right)$. By Proposition 3.11, $\sigma$-almost every $\left(x_{1}, x_{2}\right) \in X \times X$ also has the property that $\left(x_{1}, x_{2}\right) \in \operatorname{supp}\left(\lambda_{\left(x_{1}, x_{2}\right)}\right)$. Using

Lemma 3.7 it follows that $\sigma$-almost every $\left(x_{1}, x_{2}\right) \in X \times X$ satisfies $\lambda_{\left(x_{1}, x_{2}\right)}=\lambda_{\left(a, x_{1}\right)}$. We conclude that $\sigma$-almost every $\left(x_{1}, x_{2}\right)$ satisfies both of the following properties (matching Properties 1 and 2 in Section 3.1):

- $\left(a, x_{1}\right) \in \operatorname{gen}\left(\lambda_{\left(a, x_{1}\right)}, \Psi\right)$,
- $\left(x_{1}, x_{2}\right) \in \operatorname{supp}\left(\lambda_{\left(a, x_{1}\right)}\right)$.

Since orbits of generic points are dense in the support (see, eg., [16, Lemma 2.4]), we deduce that for $\sigma$-almost every $\left(x_{1}, x_{2}\right) \in X^{2}$, the point $\left(a, x_{1}, x_{2}\right) \in X^{3}$ is an Erdős progression.

To finish the proof, note that if $\left(a, x_{1}, x_{2}\right) \in X^{3}$ is an Erdős progression then $\left(a, x_{1}, x_{2}\right)$ satisfies the conclusion of Theorem 3.2, part (i), if and only if

$$
\left(x_{1}, x_{2}\right) \in E \times T^{-t} E
$$

for some $t \in \mathbb{N}$, whereas ( $a, x_{1}, x_{2}$ ) satisfies the conclusion of part (ii) if and only if

$$
\left(x_{1}, x_{2}\right) \in(T \times T)^{-t}(E \times E)
$$

for some $t \in \mathbb{N}$. Therefore, the proof is complete once we verify that

$$
\begin{equation*}
\sigma\left(E \times\left(\bigcup_{t \in \mathbb{N}} T^{-t} E\right)\right)>0 \quad \text { and } \quad \sigma\left(\bigcup_{t \in \mathbb{N}}(T \times T)^{-t}(E \times E)\right)>0 \tag{3.12}
\end{equation*}
$$

Since $\mu$ is ergodic and $E$ has positive measure, the union $\bigcup_{t \in \mathbb{N}} T^{-t} E$ covers all of $X$ up to a set of measure 0 (with respect to $\mu$ ). Writing

$$
E \times\left(\bigcup_{t \in \mathbb{N}} T^{-t} E\right)=(E \times X) \cap\left(X \times\left(\bigcup_{t \in \mathbb{N}} T^{-t} E\right)\right)
$$

we use Lemma 3.6 and then Lemma 3.5 to obtain

$$
\sigma\left(E \times\left(\bigcup_{t \in \mathbb{N}} T^{-t} E\right)\right)=\sigma(E \times X)=\mu(E)>0
$$

as desired.
We are left with showing the second positivity statement in (3.12). Write $u=$ $2 t$ when $u$ is even and $u=2 t+1$ when $u$ is odd. Since $\sigma$ is $\left(T \times T^{2}\right)$-invariant, $\sigma\left((T \times T)^{-2 t}(E \times E)\right)=\sigma\left(T^{-t} E \times E\right)$ and $\sigma\left((T \times T)^{-2 t-1}(E \times E)\right)=\sigma\left(T^{-t-1} E \times\right.$ $\left.T^{-1} E\right)$. On the other hand,

$$
\sigma\left(\bigcup_{t \in \mathbb{N}} T^{-t} E \times E \cup \bigcup_{t \in \mathbb{N}} T^{-t-1} E \times T^{-1} E\right)=\sigma\left(X \times\left(E \cup T^{-1} E\right)\right)>0
$$

Therefore, for some $u \in \mathbb{N}$ we have that $\sigma\left((T \times T)^{-u}(E \times E)\right)>0$.

## References

[1] C. Aliprantis and K. Border. Infinite dimensional analysis. A hitchhiker's guide. Third edition. Springer, Berlin, 2006. xxii+703 pp.
[2] M. Beiglböck, V. Bergelson, N. Hindman, and D. Strauss. Multiplicative structures in additively large sets. J. Combin. Theory Ser. A 113 (2006), no. 7, 1219-1242.
[3] M. Di Nasso, I. Goldbring, R. Jin, S. Leth, M. Lupini, and K. Mahlburg. On a sumset conjecture of Erdős. Canad. J. Math. 67 (2015), no. 4, 795-809.
[4] M. Einsiedler and T. Ward. Ergodic theory with a view towards number theory. Graduate Texts in Mathematics, 259. Springer-Verlag London, Ltd., London, 2011.
[5] P. Erdős and P. Turán. On some sequences of integers. ]em J. London Math. Soc. 11 (1936), 261-264.
[6] P. Erdős. Problems and results in combinatorial number theory. Astérisque 2-25 (1975), 295-309.
[7] P. Erdős. Problems and results on combinatorial number theory. III. Number theory day (Proc. Conf., Rockefeller Univ., New York, 1976), pp. 43-72. Lecture Notes in Math., Vol. 626, Springer, Berlin, 1977.
[8] P. Erdős. A survey of problems in combinatorial number theory. Ann. Discrete Math. 6 (1980), 89-115.
[9] H. Furstenberg. Ergodic behavior of diagonal measures and a theorem of Szemerédi on arithmetic progressions. J. Analyse Math. 31 (1977), 204-256.
[10] H. Furstenberg. Recurrence in Ergodic Theory and Combinatorial Number Theory. M. B. Porter Lectures. Princeton University Press, Princeton, N.J., 1981.
[11] E. Glasner, W. Huang, S. Shao, X. Ye. Regionally proximal relation of order $d$ along arithmetic progressions and nilsystems. arxiv:1911.04691
[12] N. Hindman. Finite sums from sequences within cells of a partition of $\mathbb{N}$. J. Combinatorial Theory Ser. A 17 (1974), 1-11.
[13] N. Hindman. Ultrafilters and combinatorial number theory. Number theory, Carbondale 1979 (Proc. Southern Illinois Conf., Southern Illinois Univ., Carbondale, Ill., 1979), pp. 119-184, Lecture Notes in Math., 751, Springer, Berlin, 1979.
[14] B. Host. A short proof of a conjecture of Erdős proved by Moreira, Richter and Robertson. Discrete Anal. (2019), Paper No. 19, 10 pp.
[15] B. Host and B. Kra. Nonconventional ergodic averages and nilmanifolds. Ann. of Math. (2) 161 (2005), no. 1, 397-488.
[16] B. Kra, J. Moreira, F. Richter, D. Robertson. Infinite sumsets in sets with positive density. J. Amer. Math. Soc. (2023)
[17] B. Kra, J. Moreira, F. Richter, D. Robertson. Problems on infinite sumset configurations in the integers and beyond. arXiv:2311.06197
[18] J. Moreira, F. Richter, and D. Robertson. A proof of a sumset conjecture of Erdős. Ann. of Math. (2) 189 (2019), no. 2, 605-652.
[19] M. Nathanson. Sumsets contained in infinite sets of integers. J. Combin. Theory Ser. A 28 (1980), no. 2, 150-155.
[20] K. Roth. On certain sets of integers. J. London Math. Soc. 28 (1953), 245-252.
[21] E. Szemerédi. On sets of integers containing no $k$ elements in arithmetic progression. Acta Arith. 27 (1975), 199-245.
[22] B. L. van der Waerden. Beweis einer baudetschen vermutung, Nieuw. Arch. Wisk. 15 (1927), 212-216.

Bryna Kra
Northwestern University
kra@math.northwestern.edu
Joel Moreira
University of Warwick
joel.moreira@warwick.ac.uk
Florian K. Richter
École Polytechnique Fédérale de Lausanne (EPFL)
f.richter@epfl.ch

Donald Robertson
University of Manchester
donald.robertson@manchester.ac.uk

