# THE HRT CONJECTURE FOR TWO CLASSES OF SPECIAL CONFIGURATIONS 

KASSO A. OKOUDJOU AND VIGNON OUSSA


#### Abstract

The HRT (Heil-Ramanathan-Topiwala) conjecture stipulates that the set of any finitely many time-frequency shifts of a non-zero square Lebesgue integrable function is linearly independent. The present work settles two special cases of this conjecture, namely, the cases where the set of time-frequency shifts has cardinality $N+1$ such that either $N$ of the points lie on some integer lattice and the last point is arbitrary, or $N$ of the points are on a line, while the last point does not belong this line. In both cases, we prove that the HRT conjecture holds appealing mainly to various forms of the ergodic theorem. We note that, in recent years, the latter case has been the subject of many investigations-notably, the subcase where $N=3$ - and our work completely resolves it.


## 1. Introduction

For a nonzero function $f \in L^{2}(\mathbb{R})$, a time-frequency shift of $f$ is a function of the type

$$
M_{y} T_{x} f(t)=e^{-2 \pi i y t} f(t-x) \quad(x, y) \in \mathbb{R}^{2}
$$

The operator $M_{y}$ is called a modulation or frequency-shift operator, and $T_{x}$ is an operator acting by shifting $f$ in time domain. Generally, modulation and translation operators do not commute, and in fact, these two families of unitary operators generate a non-commutative group called the Heisenberg group. The ubiquitous nature of the Heisenberg group and its relevance across numerous subjects within harmonic analysis is a fact that has been extensively documented in the literature [17]. For instance, the Heisenberg group is fundamental to the foundation of time-frequency or Gabor analysis and is a source of important examples in frame theory [12, 15]. For a historical development, as well as an in-depth treatment of Gabor analysis and connection with topics such as frame theory, we refer the reader to [7, 8, 12, 14].

We recall that given a countable subset $\Lambda$ of $\mathbb{R}^{2}$, the collection of functions

$$
\mathcal{G}(f, \Lambda)=\left\{M_{y} T_{x} f:(x, y) \in \Lambda\right\}
$$

is called a Gabor system. For comparison, the system obtained by substituting modulation for dilation operators is often called a wavelet system. The underlying groups for these
systems: the Heisenberg group and the affine group, share many striking similarities, as documented in [23, 22]. However, there are also several instances in which these systems behave in drastically different ways. For example, it is well-known that there exist finite systems produced by time-dilation operators which are linearly dependent. This property guarantees the existence of a scaling function as a solution to a suitable functional equation-a central ingredient in the construction of multiresolution orthonormal wavelets [15]. However, to this date, we do not know the extent to which an analogous statement holds for finite systems of time-frequency shifts. In fact, all progress made in this direction suggests that any finite system of time-frequency shifted copies of a nonzero function is linearly independent. This conjecture, first formulated by Heil, Ramanathan, and Topiwala, is now known as the HRT Conjecture [16] and may be stated as follows.

Conjecture 1. [16] (HRT Conjecture) Given a fixed finite set $\Lambda \subset \mathbb{R}^{2}$, and a nonzero function $f \in L^{2}(\mathbb{R})$, the associated finite Gabor system, $\mathcal{G}(f, \Lambda)$ is linearly independent.

For the specific case where all the points lie on the same line (which is equivalent to stating that all modulation parameters are identical or all the translation parameters are identical), a straightforward application of the Fourier transform verifies the conjecture. Generally, however, the problem is much more challenging than one might initially anticipate. To this date, the most general result is attributed to Linnell [18], who was able to corroborate the conjecture for the case where $\Lambda$ is a shift of a finite subset of a discrete subgroup of $\mathbb{R}^{2 d}, d \in \mathbb{N}$.

In terms of the number of points, the smallest case not handled by Linnell's result can only occur for specific configurations of a set $\Lambda$ containing four points. To date, the status of the HRT conjecture for sets of four points can be summarized as follows.

Proposition 1. Let $\Lambda \subset \mathbb{R}^{2}$ be such that $\# \Lambda=4$, and let $0 \neq f \in L^{2}(\mathbb{R})$. The finite Gabor system $\mathcal{G}(f, \Lambda)$ is linearly independent in each of the following cases:
(1) $\Lambda$ is $(2,2)$-configuration, that is two of the points are on a line and the other two on a parallel line, and $0 \neq f \in L^{2}(\mathbb{R})$ is arbitrary, [4, 5];
(2) $\Lambda$ is 13-configuration, that is three of the points are on a line and the fourth point off that line, and $0 \neq f \in \mathcal{S}(\mathbb{R})$ [4];
(3) $\Lambda$ is an arbitrary set of four points, but extra restrictions are imposed on $f$, e.g., [2, 3, 21].

Still, the general problem of the HRT conjecture for four points remains unsettled, and the main aim of this paper is to present new results, including the proof of the HRT conjecture for all $(1,3)$ configurations and all $L^{2}$ functions.

Our work is motivated by the two special cases of 4-point HRT, namely the $(1,3)$ configurations, as well as well the case in which three of the points are on the $\mathbb{Z}^{2}$-lattice and the fourth is arbitrary.


Figure 1. Examples of four-point configurations considered in the sequel.
For example, the configuration $\left\{T_{x_{1}} f, T_{x_{2}} f, T_{x_{3}} f, M f\right\}$, for arbitrary $x_{1}, x_{2}, x_{3} \in \mathbb{R}$ falls into the first category, while the configuration $\left\{f, M f, T f, M_{\alpha} T_{\beta} f\right\}$ falls into the second, see Figure 1 .

Our results will deal with generalization of these two cases to the setting of $N+1$ timefrequency shifts with $N \geq 3$. In particular, our results are summarized as follows.

Theorem 2. Let $\left\{\left(x_{k}, y_{k}\right)\right\}_{k=1}^{N} \cup\{(\alpha, \beta)\} \subset \mathbb{R}^{2}$ be such that, either
(i) $\left\{\left(x_{k}, y_{k}\right)\right\}_{k=1}^{N} \subset \mathbb{Z}^{2}$, and $(\alpha, \beta) \in \mathbb{R}^{2} \backslash \mathbb{Z}^{2}$, or
(ii) $\left\{\left(x_{k}, y_{k}\right)\right\}_{k=1}^{N} \subset \mathbb{R}^{2}$ lie on a line $L$ and $(\alpha, \beta) \in \mathbb{R}^{2} \backslash L$.

If $f \in L^{2}(\mathbb{R})$ is such that for some nonzero complex coefficients, $c_{1}, \cdots, c_{N}$ the following linear dependency relation holds:

$$
\sum_{k=1}^{N} c_{k} M_{y_{k}} T_{x_{k}} f=M_{\alpha} T_{\beta} f
$$

then $f$ must be identically zero.
As will be seen in most of our proofs, Ergodic theory plays a critical role. This is primarily due to the structure of the configurations of interest in this work. For any time-frequency equations of the type $\sum_{k=1}^{N} c_{k} M_{y_{k}} T_{x_{k}} f=M_{\alpha} T_{\beta} f$, or $\sum_{k=1}^{N} c_{k} T_{x_{k}} f=M f$, we prove that there exist unitary operators that diagonalize the left hand side of each equation. This allows us to reformulate each time-frequency equation in the form $|p(z) F(z)|=|F(z+\gamma)|$ with a varying parameter $z$, some polynomial $p$ in the variable $z$, and some fixed parameter $\gamma$. The overarching strategy consists of step-wise iterations of $|p(z) F(z)|=|F(z+\gamma)|$ combined with suitable applications of various versions of the Ergodic Theorem and its relatives (such
as the Poincaré Recurrence Theorem.) For each of these applications, we refer the interested reader to the following relevant textbooks [6, 24] for detailed treatment.

The rest of the paper is organized as follows. In Section 2, we consider the case of $\left(1, \mathbb{Z}^{2}\right)$ configurations, and state and prove our first main result, namely Theorem 3. In Section 3, we prove our second main result, Theorem 6, which settles the HRT conjecture for all $(1, N)$ configurations.

## 2. The HRT conjecture for $\left(1, \mathbb{Z}^{2}\right)$-Configurations

A collection of $N+1$ points $\Lambda=\left\{\left(x_{k}, y_{k}\right)\right\}_{k=1}^{N} \cup\{(\alpha, \beta)\} \subset \mathbb{R}^{2}$ is a $\left(1, \mathbb{Z}^{2}\right)$-configuration if $\left\{\left(x_{k}, y_{k}\right)\right\}_{k=1}^{N} \subset \mathbb{Z}^{2}$ and $(\alpha, \beta)$ is arbitrary. We refer to Figure 1 for an illustration in the case $N=3$. Our first result establishes the HRT conjecture for all such configurations.

Theorem 3. Fix $\left\{\left(x_{k}, y_{k}\right): 1 \leq k \leq N\right\} \subset \mathbb{Z}^{2}$, and let $(\alpha, \beta) \in \mathbb{R}^{2} \backslash \mathbb{Z}^{2}$. If there exists a function $f \in L^{2}(\mathbb{R})$ such that

$$
\begin{equation*}
\sum_{k=1}^{N} c_{k} M_{y_{k}} T_{x_{k}} f=M_{\alpha} T_{\beta} f \tag{1}
\end{equation*}
$$

for some nonzero complex coefficients $c_{1}, \cdots, c_{N}$, then $f$ must be the zero function.
An immediate consequence of Theorem 3, is the follwing result which answers a subconjecture of C. Heil, [15, Conjecture 11.39].

Corollary 4. For any function $0 \neq f \in L^{2}(\mathbb{R})$, the finite collection of time-frequency shifts $\left\{f, M_{1} f, T_{1} f, M_{\alpha} T_{\beta} f\right\}$ is linearly independent in $L^{2}(\mathbb{R})$ for any point $(\alpha, \beta) \in \mathbb{R}^{2}$.

Theorem 3 can also be extended to ( $1, A \mathbb{Z}^{2}$ )-configurations where $A$ is any $2 \times 2$ symplectic matrix, that is an invertible matrix whose determinant is 1 .

Corollary 5. Suppose that $A$ is a $2 \times 2$ matrix with $\operatorname{det} A=1$. The HRT conjecture is true for any $0 \neq f \in L^{2}(\mathbb{R})$ and any set of $N+1$ point of the form $\left\{\left(x_{k}, y_{k}\right): 1 \leq k \leq N\right\} \cup\{(\alpha, \beta)\}$ where $\left\{\left(x_{k}, y_{k}\right): 1 \leq k \leq N\right\} \subset A \mathbb{Z}^{2}$, and $(\alpha, \beta) \in \mathbb{R}^{2}$.

Proof. When $(\alpha, \beta) \in A \mathbb{Z}^{2}$, we are back in the lattice case, which was already settled by Linell. Therefore, we may assume that $(\alpha, \beta) \in \mathbb{R}^{2} \backslash A \mathbb{Z}^{2}$, and suppose that there, HRT was false for a function $f$ and one such set of $N+1$ points. We may then use the fact that $A$ is a symplectic matrix to find a function $g$ (image of $f$ under a metaplectic transformation) such that $g$ and the set of points $\left\{\left(x_{k}^{\prime}, y_{k}^{\prime}\right): 1 \leq k \leq N\right\} \cup\left\{\left(\alpha^{\prime}, \beta^{\prime}\right)\right\}$ where $\left\{\left(x_{k}^{\prime}, y_{k}^{\prime}\right): 1 \leq k \leq N\right\} \subset$ $\mathbb{Z}^{2}$. But this will contradict Theorem 3.
2.1. Discussion of the proof of Theorem 3. As we will see, the concept of convergence and divergence of infinite products will also be of central importance in our proofs. As such, we would like to review here some of the important notions needed in the sequel, and we
refer to [1] for details. Let $\left(a_{j}\right)_{j \geq 1}$ be a sequence of nonzero real or complex numbers. If the sequence of partial products

$$
\left(\prod_{j=1}^{n} a_{j}\right)_{n \in \mathbb{N}}
$$

converges to a nonzero limit then we say that the infinite product $\prod_{j=1}^{\infty} a_{j}$ is convergent. Suppose a finite number of factors $a_{j}$ are equal to zero, and the infinite product obtained by removing these factors converges. In that case, we say that the infinite product converges to zero. Otherwise, we define the infinite product as divergent to zero. We note that if $\prod_{j=1}^{\infty} a_{j}$ is convergent, then $\lim _{j \rightarrow \infty} a_{j}=1$, see [1].

To prove Theorem 3, we note that any finite set of time-frequency operators parametrized by a collection of finite points on an integer lattice in $\mathbb{R}^{2}$, forms a collection of pairwise commuting operators. As such, the Spectral Theorem [13] guarantees the existence of a unitary operator which diagonalizes the operator

$$
J=\sum_{k=1}^{N} c_{k} M_{y_{k}} T_{x_{k}}
$$

In our case, this unitary operator is the Zak transform $Z: L^{2}(\mathbb{R}) \rightarrow L^{2}\left([0,1]^{2}\right)$ defined formally, by

$$
\begin{equation*}
Z f(t, \omega)=\sum_{k \in \mathbb{Z}} f(t+k) e^{-2 \pi k \omega i} \tag{2}
\end{equation*}
$$

We summarize below the main properties of the Zak transform and refer to [12] for a complete introduction to this transform.

Lemma 1. The Zak transform $Z$ satisfies the following properties.
(1) $Z$ is a unitary operator $L^{2}(\mathbb{R})$ onto $L^{2}\left([0,1)^{2}\right)$.
(2) For $\varphi \in L^{2}(\mathbb{R})$ and $(t, \omega) \in[0,1)^{2}$, we have:
(a) $[Z T \varphi](t, \omega)=e^{-2 \pi \omega i} \cdot[Z \varphi](t, \omega)$.
(b) $[Z M \varphi](t, \omega)=e^{-2 \pi i t} \cdot[Z \varphi](t, \omega)$.
(c) For each $\alpha, \beta>0,\left[Z M_{\alpha} T_{\beta} \varphi\right](t, \omega)=e^{-2 \pi i \alpha t} \cdot[Z \varphi](t-\beta, \omega+\alpha)$.
(d) For any integer $j,[Z T \varphi](t+j, \omega)=e^{2 \pi \omega j i}[Z T \varphi](t, \omega)$.
(e) For any integer $j,[Z T \varphi](t, \omega+j)=[Z T \varphi](t, \omega)$.

Applying Zak transform to the linear dependence equation $J f=M_{\alpha} T_{\beta} f$, and letting $F=Z f$, for almost every $(t, w) \in[0,1]^{2}$, we have

$$
\begin{equation*}
\sum_{k=1}^{N} c_{k} e^{-2 \pi i y_{k} t} e^{-2 \pi i x_{k} \omega} \cdot F(t, \omega)=e^{-2 \pi i \alpha t} \cdot F(t-\beta, \omega+\alpha) \tag{3}
\end{equation*}
$$

or equivalently,

$$
\begin{equation*}
p(t, \omega) \cdot F(t, \omega)=e^{-2 \pi i \alpha t} \cdot F(t-\beta, \omega+\alpha) \tag{4}
\end{equation*}
$$

where $p: \mathbb{R}^{2} \rightarrow \mathbb{C}$ is the non-zero trigonometric polynomial given by

$$
\begin{equation*}
p(t, \omega)=\sum_{k=1}^{N} c_{k} e^{-2 \pi i y_{k} t} e^{-2 \pi i x_{k} \omega}=\sum_{k=1}^{N} c_{k} e^{-2 \pi i\left\langle(t, \omega),\left(y_{k}, x_{k}\right)\right\rangle} . \tag{5}
\end{equation*}
$$

It follows that for almost every $(t, w) \in[0,1]^{2}$,

$$
|p(t, \omega)| \cdot|F(t, \omega)|=|F(t-\beta, \omega+\alpha)| \Longleftrightarrow|F(z+\gamma)|=|p(z)| \cdot|F(z)|
$$

where, for simplicity in the notations, we set $\gamma=(-\beta, \alpha)$ and $z=(t, \omega)$. Iterating, the last equation, we see that for almost every $z \in[0,1]^{2}$ and for any positive integer $n \geq 1$,

$$
\begin{equation*}
|F(z+n \gamma)|=\left|\prod_{j=0}^{n-1} p(z+j \gamma)\right| \cdot|F(z)| \tag{6}
\end{equation*}
$$

We let $S_{F}$ be the essential support of $F$ in $[0,1]^{2}$ and let $K_{F}=[0,1]^{2} \backslash S_{F}$ be its complement. By Lusin's theorem, there exists a set of positive measure $\Lambda \subset S_{F}$, such that the restriction $|F|_{\mid \Lambda}$ of $|F|$ onto $\Lambda$ is continuous. Finally, owning to the analyticity of $|p(z)|^{2}$, we may assume that $\Lambda$ contains no zeroes of $p(z)$.

The final tool needed in our proof is the measure-preserving map $\sigma:[0,1)^{2} \rightarrow[0,1)^{2}$ such that for any $z \in[0,1)^{2}, \sigma(z)$ is the unique element in $[0,1)^{2}$ which is a solution to the equation $z+\gamma \mathbb{Z}^{2}=\sigma(z)+\mathbb{Z}^{2}$. In other words, from here on, we will regard the range of $\sigma$ as a connected transversal for the $\mathbb{Z}^{2}$-cosets in $\mathbb{R}^{2}$.

Observe that, (6) suggests that, at least for continuous functions, the behavior of $|F|$ at $z+n \gamma$ is fully determined by its value at a single point such as $z$. Thus, if for instance, $F$ is assumed to be continuous, since $p$ is predetermined; for any unknown function $f$ for which the HRT conjecture is true, the restriction of $|F|$ to the set $\{z+n \gamma: n \in \mathbb{N}\}$ is completely determined by its value at a single element such as $z$. Even if, $F$ is not assumed to be continuous, since the Zak transform converges almost everywhere and since (6) imposes a rather stringent condition on $F$, and, in principle, $J f=M_{\alpha} T_{\beta} f$ should be incompatible with the square-integrability of nonzero $f$. Our main concern in this section is to fully substantiate these claims.

### 2.2. Proof of Theorem 3.

Proof of Theorem [3. The proof is twofold. First, we assume that $\{1, \alpha, \beta\}$ is linearly independent over the rationals. Secondly, we deal with the situation where this set is dependent over the rationals.

Case 1: Let us assume that $\{1, \alpha, \beta\}$ is linearly independent over the rationals. It follows that for any $z \in[0,1]^{2}$, the set $\{\sigma(z+n \gamma): n \in \mathbb{N}\}$ is dense in $[0,1]^{2}$, and that $\sigma$ defines an ergodic map on the unit square (viewed as a fundamental domain of an integer lattice). Furthermore, according to Poincaré Recurrence Theorem, for almost every $\lambda \in \Lambda, \sigma^{n}(\lambda) \in \Lambda$ for infinitely many $n$.

Next, assume that that both $S_{F}$ (the support of $F$ ) and its complement $K_{F}$ in $S$, have positive Lebesgue measure as subsets of the unit square. It follows then, that there exists $n \in$ $\mathbb{N}$, such that $\sigma^{n}\left(S_{F}\right) \cap K_{F}$ has positive Lebesgue measure in $[0,1]^{2}$, which is a contradiction (see Equation (6)). Since $F \neq 0$, we conclude that $\left|S_{F}\right|=1$ (or that $\left|K_{F}\right|=0$ ), where $|E|$ denotes the Lebesgue measure of $E \subset \mathbb{R}^{2}$.

We now see that $n \geq 1$ and any $z \in S_{F}$, (6) is equivalent to

$$
\begin{equation*}
\left|\frac{F\left(\sigma^{n}(z)\right)}{F(z)}\right|=\left|\prod_{j=0}^{n-1} p\left(\sigma^{j}(z)\right)\right| . \tag{7}
\end{equation*}
$$

Let $z_{0} \in S_{F}$ and fix $\lambda_{0} \in \Lambda$. Choose a subsequence $\left\{\sigma^{n_{k}}\left(\lambda_{0}\right)\right\}_{k \geq 1} \subset \Lambda$ such that $\lim _{k \rightarrow \infty} \sigma^{n_{k}}\left(\lambda_{0}\right)=$ $z_{0}$. Appealing to (7), we derive the following:

$$
\left|\frac{F\left(z_{0}\right)}{F\left(\lambda_{0}\right)}\right|=\lim _{k \rightarrow \infty}\left|\frac{F\left(\sigma^{n_{k}}\left(\lambda_{0}\right)\right)}{F\left(\lambda_{0}\right)}\right|=\lim _{k \rightarrow \infty}\left|\prod_{j=0}^{n_{k}-1} p\left(\sigma^{j}\left(\lambda_{0}\right)\right)\right| .
$$

Moreover, applying the natural $\log$ to each side of the previous equation, the following is immediate:

$$
\lim _{k \rightarrow \infty} \sum_{j=0}^{n_{k}-1} \ln \left|p\left(\sigma^{j}\left(\lambda_{0}\right)\right)\right|=\ln \left|\frac{F\left(z_{0}\right)}{F\left(\lambda_{0}\right)}\right|<\infty .
$$

Set $s_{k}=\sum_{j=0}^{n_{k}-1} \ln \left|p\left(\sigma^{j}\left(\lambda_{0}\right)\right)\right|$ and note that $\lim _{k \rightarrow \infty} s_{k+1}-s_{k}=0$, which implies that

$$
\sum_{k=1}^{\infty}\left(s_{k+1}-s_{k}\right)=\lim _{K \rightarrow \infty} \sum_{k=1}^{K}\left(s_{k+1}-s_{k}\right)=\ln \left|\frac{F\left(z_{0}\right)}{F\left(\lambda_{0}\right)}\right|-s_{1}<\infty
$$

Now the sequence $\left\{s_{k+1}-s_{k}\right\}_{k \geq 1}$ can be written as $\left\{\ln \left|p\left(\sigma^{m_{k}}\left(\lambda_{0}\right)\right)\right|\right\}_{k \geq 1}$ where the increasing sequence $\left\{m_{k}\right\}$ is the enumeration of $\left\{n_{k}+\ell: \ell=0,1,2, \ldots, n_{k+1}-1-n_{k}, k=1,2, \ldots\right\}$. We observe that the original sequence $\left\{\ln \left|p\left(\sigma^{n_{k}}\left(\lambda_{0}\right)\right)\right|\right\}_{k \geq 1}$ is a subsequence of the new sequence
$\left\{\ln \left|p\left(\sigma^{m_{k}}\left(\lambda_{0}\right)\right)\right|\right\}_{k \geq 1}$. We can then write that

$$
\sum_{k=1}^{\infty} \ln \left|p\left(\sigma^{m_{k}}\left(\lambda_{0}\right)\right)\right|<\infty
$$

This implies that

$$
\lim _{k \rightarrow \infty} \ln \left|p\left(\sigma^{m_{k}}\left(\lambda_{0}\right)\right)\right|=0 \Longrightarrow \lim _{k \rightarrow \infty} \ln \left|p\left(\sigma^{n_{k}}\left(\lambda_{0}\right)\right)\right|=0
$$

and as such, $\lim _{k \rightarrow \infty}\left|p\left(\sigma^{n_{k}}\left(\lambda_{0}\right)\right)\right|=\left|p\left(z_{0}\right)\right|=1$.
We conclude that $|p(z)|=1$ on $[0,1]^{2}$. Note however, that $|p|$ is not a constant function, since $|p|^{2}$ is real-analytic (a trigonometric polynomial, more precisely.) This produces the desired contradiction.

Case 2: Assume now that the set $\{1, \alpha, \beta\}$ is rationally dependent. Since the case where the parameters $\alpha, \beta$ are rationals (see [19, Proposition 6.3] and [18]) has already been settled, to complete the proof, it remains to address the subcase where either $\alpha$ or $\beta$ is irrational. Without loss of generality, assume that $0 \neq \beta \notin \mathbb{Q}$ and write $\alpha=r+s \beta$ for some rational numbers $r, s$. For each $\lambda \in[0,1)^{2}$, let

$$
L(\lambda)=\{\lambda+t \gamma \quad \bmod 1: t \in \mathbb{R}\}
$$

be the toral line passing through $\lambda$ and $\gamma$. Then $L(\lambda)$ can be identified with a finite union of parallel line segments inside the two-dimensional unit square. Note that these toral lines partition the unit square $[0,1)^{2}$.


Figure 2. An illustration of several distinct toral lines in $[0,1)^{2}$

In this case, $\sigma$ fails to act ergodically on the (entire) unit square. Next, we aim to show that the trigonometric polynomial $p$ in (5) is invariant under the action of $\sigma$. Indeed, for $t \in \mathbb{R}$, let $z_{t} \in[0,1)^{2}$ be the unique element such that $\lambda+t \gamma-z_{t}=q \in \mathbb{Z}^{2}$. It follows that $\sigma(\lambda+t \gamma)=z_{t}$, and that

$$
\begin{aligned}
p(\sigma(\lambda+t \gamma)) & =\sum_{k=1}^{N} c_{k} e^{-2 \pi i\left\langle\sigma(\lambda+t \gamma),\left(y_{k}, x_{k}\right)\right\rangle}=\sum_{k=1}^{N} c_{k} e^{-2 \pi i\left\langle z_{t},\left(y_{k}, x_{k}\right)\right\rangle} \\
& =\sum_{k=1}^{N} c_{k} e^{-2 \pi i\left\langle z_{t}+q,\left(y_{k}, x_{k}\right)\right\rangle}=p(\lambda+t \gamma)
\end{aligned}
$$

By the continuity of $p$, we conclude that $p(\sigma(z))=p(z)$ for each $z \in L(\lambda)$. In addition, since the action of $\sigma$ is ergodic (with respect to a normalized arclength measure) on the toral lines $L(\lambda)$, it follows that $p$ must be constant on each of these lines.

We proceed as in the first case and choose a set of positive measure $\Lambda \subset S_{F}$ such that the restriction $|F|_{\mid \Lambda}$ of $|F|$ onto $\Lambda$ is continuous. Let $\lambda \in \Lambda$. Using the fact that $\sigma$ is measure-preserving, once again, we may apply Poincaré Recurrence Theorem to assert the existence of a subsequence $\left(\sigma^{n_{k}}(\lambda)\right)_{k \in \mathbb{N}} \subset \Lambda$ such that $\lim _{k \rightarrow \infty} \sigma^{n_{k}}(\lambda)=\lambda_{0} \in \Lambda$. Using the same arguments as in Case 1, we conclude that

$$
\left|\frac{F\left(\lambda_{0}\right)}{F(\lambda)}\right|=\lim _{k \rightarrow \infty}\left|\frac{F\left(\sigma^{n_{k}}(\lambda)\right)}{F(\lambda)}\right|=\lim _{k \rightarrow \infty}\left|\prod_{j=0}^{n_{k}-1} p\left(\sigma^{j}(\lambda)\right)\right|,
$$

which implies that $\lim _{k \rightarrow \infty}\left|p\left(\sigma^{n_{k}}(\lambda)\right)\right|=\left|p\left(\lambda_{0}\right)\right|=1$. We now use the fact that $\left(\sigma^{n_{k}}(\lambda)\right)_{k \in \mathbb{N}} \subset$ $L(\lambda)$ to conclude that $\lambda_{0} \in L(\lambda)$. It follows that $|p(\lambda)|=\left|p\left(\lambda_{0}\right)\right|=1$.

Since the collection of toral lines $L(\lambda)$ is a partition of the unit square (up to a set of Lebesgue measure zero) and since $\Lambda$ is assumed to have positive Lebesgue measure, then $|p|^{2}=1$ on a set of positive Lebesgue measure. Hence $|p|^{2}$ (being real analytic), and $|p|$ must be constant everywhere on the unit square. This would imply that all but one of the coefficients $c_{k} s$ is nonzero. This produces a desired contradiction.

## 3. The HRT for $(1, n)$-configurations

In this section, we prove the HRT conjecture for all $(1, n)$-configurations. In particular, the result settles the $(1,3)$-configurations and should be compared to the results in [4, 5, 20, 21].

Theorem 6. Suppose that $\left\{x_{k}\right\}_{k=1}^{N} \subset \mathbb{R}$ is a set of $N$ distinct points. Assume that there exists a function $f \in L^{2}(\mathbb{R})$ such that

$$
\begin{equation*}
\sum_{k=1}^{N} c_{k} T_{x_{k}} f=M f \tag{8}
\end{equation*}
$$

for some nonzero complex coefficients $c_{1}, \cdots, c_{N}$. Then $f$ must be the zero function.
3.1. Discussion of the proof of Theorem 6. Since the case where all translation parameters $x_{k}$ are rational is already established [19, 18, Proposition 6.3], without loss of generality, we may assume that

$$
\operatorname{dim}_{\mathbb{Q}}\left(\mathbb{Q}+\mathbb{Q} x_{1}+\cdots+\mathbb{Q} x_{N}\right)>1
$$

Although the rationality condition above implies that at least, one of the translation parameters must be irrational, it is perhaps worth pointing out that it is not necessarily the case that $\left\{1, x_{1}, \cdots, x_{N}\right\}$ is linearly independent over the rationals. Applying Fourier transform to each side of Equation (8) and letting $F$ be the Fourier transform of $f$, it follows that

$$
\begin{equation*}
\left(\sum_{k=1}^{N} c_{k} e^{-2 \pi i x_{k} \xi}\right) \cdot F(\xi)=F(\xi+1) \text { a.e. } \tag{9}
\end{equation*}
$$

Next, we consider the trigonometric polynomial

$$
p(\xi)=\sum_{k=1}^{N} c_{k} e^{-2 \pi i x_{k} \xi}
$$

The following will be quite convenient for the proof of our main result.
Lemma 2. $p$ is a matrix coefficient of an $N$-dimensional unitary representation $\mathbb{R}$.
Proof. Indeed, it is not hard to verify that $p(\xi)=\left\langle e^{\xi D} c, v\right\rangle$ where

$$
D=\left[\begin{array}{lll}
-2 \pi i x_{1} & & \\
& \ddots & \\
& & -2 \pi i x_{N}
\end{array}\right] e^{\xi D}=\left[\begin{array}{ccc}
e^{-2 \pi i \xi x_{1}} & & \\
& \ddots & \\
& & e^{-2 \pi i \xi x_{N}}
\end{array}\right]:=A
$$

Additionally

$$
v=\left[\begin{array}{c}
1 \\
\vdots \\
1
\end{array}\right], c=\left[\begin{array}{c}
c_{1} \\
\vdots \\
c_{N}
\end{array}\right] .
$$

Each quantity $p(\xi)=\left\langle e^{\xi D} c, v\right\rangle$ is obtained by letting $\mathbb{R}$ acts on $\mathbb{C}^{N}$ via the action $(\xi, c) \mapsto e^{\xi D} c$ and this produces a family of scalars $\left(\left\langle e^{\xi D} c, v\right\rangle\right)_{\xi \in \mathbb{R}}$ obtained by projecting $e^{\xi D} c$ continuously on the function $v$ described above. As $\xi$ runs through the set of real numbers, $e^{\xi D} c$ traces out a curve wrapping itself densely around a torus in $\mathbb{C}^{N}$ passing through, c.

Observe, that generally, the countable group generated by matrix $A=e^{D}$ is not closed in the subspace topology of the General Linear Group which $A$ belongs to. In what follows, let
$G$ be the topological closure of the countable group

$$
\Gamma=\left\{A^{j}=e^{j D}: j \in \mathbb{Z}\right\}
$$

generated by $A$. Additionally, let $H$ be the topological closure (as a subset of the General Linear Group) of the 1-parameter group $\left\{e^{\xi D}: \xi \in \mathbb{R}\right\}$. We refer to [11] for more background.

Lemma 3. $G$ is a compact abelian group which is either homeomorphic to a torus or to a finite union of translates of tori of the same positive dimension.

Proof. By assumption, since $\operatorname{dim}_{\mathbb{Q}}\left(\mathbb{Q}+\mathbb{Q} x_{1}+\cdots+\mathbb{Q} x_{N}\right)>1$, the set $\left\{x_{1}, \cdots, x_{N}\right\}$ must contain at least one irrational number. Without loss of generality, we may also assume that $x_{1} \notin \mathbb{Q}$, and additionally, $\left\{1, x_{1}, \cdots, x_{m}\right\}$ is a basis for the rational space $\mathbb{Q} x_{1}+\cdots+\mathbb{Q} x_{m}+\mathbb{Q}$. In other words, we may rewrite the remaining parameters $x_{m+1}, \cdots, x_{N}$ as follows:

$$
\begin{aligned}
x_{m+1} & =u_{m+1}+\sum_{l=1}^{m} d_{m+1, l} x_{l} \text { for }\left(u_{m+1}, d_{m+1,1, \ldots} d_{m+1, m}\right) \in \mathbb{Q}^{m+1} \\
& \vdots \\
x_{N} & =u_{N}+\sum_{l=1}^{m} d_{N, l} x_{l} \text { for }\left(u_{N}, d_{N, 1, \ldots} d_{N, m}\right) \in \mathbb{Q}^{m+1} .
\end{aligned}
$$

Furthermore, we also assume that every member of the set $\left\{1, x_{1}, \cdots, x_{N}\right\}$ has a rationalized denominator. Next, let $L$ be the least common multiple of all of these denominators. For every integer $j$, let $j=q L+r$, that is $r=j \bmod L$. It follows that the diagonal entries of $A^{j}$ where $A=e^{D}$ are
$\left(e^{-2 \pi i x_{1}(q L+r)}, \cdots, e^{-2 \pi i x_{m}(q L+r)}, e^{-2 \pi i\left(u_{m+1}+\sum_{l=1}^{m} d_{m+1, l} x_{l}\right)(q L+r)}, \cdots, e^{-2 \pi i\left(u_{N}+\sum_{l=1}^{m} d_{N, l} x_{l}\right)(q L+r)}\right)$. Letting $z_{1}=e^{-2 \pi i x_{1}}, \cdots, z_{m}=e^{-2 \pi i x_{m}}$, we obtain $A^{j}=A^{q L+r}=A^{r} D_{(u, q, L)}$ where $D_{(u, q, L)}$ is a diagonal matrix with diagonal entries

$$
z_{1}^{q L}, \cdots, z_{m}^{q L}, e^{-2 \pi i\left(u_{m+1}\right)(q L)}\left(\prod_{l=1}^{m} z_{l}^{d_{m+1, l} q L}\right), \cdots, e^{-2 \pi i\left(u_{N}\right)(q L)}\left(\prod_{l=1}^{m} z_{l}^{d_{N, l} q L}\right)
$$

In summary, the closure of the countable group generated by $A$ is a topological closed subgroup of the General Linear Group which is also of the form

$$
G=\bigcup_{r=0}^{L-1} A^{r} K
$$

where $K$ is a connected component of the identity of the group.
Remark 7. In addition to Lemma 3, we also need the following facts.
(1) Since $G$ is a locally compact group, up to normalization, $G$ admits a unique measure which is invariant under left translations by elements of $G$. Moreover, this invariant
measure on $G$ is the sum of measures $\frac{1}{L} \sum_{k=1}^{L} \mu_{k}$ where each $\mu_{k}$ is the push-forward of the invariant measure $\mu_{0}$ on the connected component $K$ to $A^{k} K$ via the translation map $g \mapsto A^{k} g$.
(2) Generally $H$ and $G$ do not coincide and since $G$ is a closed subgroup of $H$, there exists a subset $Z$ of $H$ such that $\{z G: z \in Z\}$ is a transversal for $H / G$ [9]. Additionally, the transversal $Z$ admits up to normalization, a unique $H$-invariant measure. We endow each set $z G$ with the following measure. For a suitable subset $z E$ of $z G$, $\mu_{z G}(z E)=\mu_{G}(E)$ where $\mu_{G}$ is a fixed invariant measure on $G$. Additionally, we shall assume that $\mu_{G}$ is normalized to satisfy $\mu_{G}(G)=1$.
(3) Note that the group generated by $e^{D}$ is generally, not a dense subset of $H$. Therefore, rotation by $e^{D}$ is generally not an ergodic map on $H$. However, $e^{D}$ does act ergodically on every coset of $G$ in $H$. Precisely, for each fixed $z$ in a transversal $Z$, of $H / G, G$ acts on $z G$ as follows. Given $g \in G$ and $z h \in z G, g z h=z g h=z(g h) \in z G$. Since the countable subgroup $\left\{e^{j D}: j \in \mathbb{Z}\right\}$ is dense in $G$, for every fixed $z \in Z$ where $Z \subseteq H$ is a measurable transversal for $H / G$, the map $\sigma: z G \rightarrow z G$ defined by $\sigma(z h)=e^{D} z h=z e^{D} h$ is an ergodic map with respect to the measure $\mu_{z G}$.

We are now ready to prove Theorem 6.

### 3.2. Proof of Theorem 6.

Proof. Let $S$ be the support of $F$ and assume that $S$ does not contain any zero of the polynomial $p$ (this is possible since the set of these zeroes is at most countable.) For each $\xi \in S$, we can iterate (9) so that for any positive integer $n$, the following holds true:

$$
\begin{equation*}
|F(\xi+n)|=\left(\prod_{j=0}^{n-1}|p(\xi+j)|\right)|F(\xi)| \text { a.e. } \tag{10}
\end{equation*}
$$

and additionally

$$
\begin{equation*}
|F(\xi-n)|=\left(\prod_{j=1}^{n}|p(\xi-j)|\right)^{-1}|F(\xi)| \text { a.e. } \tag{11}
\end{equation*}
$$

Since $F$ is square-integrable,

$$
\infty>\int_{S}|F(\xi)|^{2} d \xi=\int_{S \cap[0,1]} \sum_{k=-\infty}^{\infty}|F(\xi+k)|^{2} d \xi
$$

This implies that for a.e. $x \in S \cap[0,1]$, we have

$$
\sum_{k=-\infty}^{\infty}|F(\xi+k)|^{2}=|F(\xi)|^{2}+\sum_{k=1}^{\infty}\left(|F(\xi+k)|^{2}+|F(\xi-k)|^{2}\right)<\infty
$$

Appealing to (10) and (11), we derive the following

$$
\begin{equation*}
\sum_{k=-\infty}^{\infty}|F(\xi+k)|^{2}=|F(\xi)|^{2} \cdot\left(1+\sum_{k=1}^{\infty}\left(\prod_{j=0}^{k-1}|p(\xi+j)|^{2}+\prod_{j=1}^{k}|p(\xi-j)|^{-2}\right)\right)<\infty \tag{12}
\end{equation*}
$$

By Egorov's theorem, there exists a set $\Lambda \subset S \cap[0,1]$ of positive Lebesgue measure such that

$$
\lim _{k \rightarrow \infty} \prod_{j=0}^{k-1}|p(\xi+j)|^{2}=\prod_{j=0}^{\infty}|p(\xi+j)|^{2}=\lim _{k \rightarrow \infty} \prod_{j=1}^{k}|p(\xi-j)|^{2}=\prod_{j=1}^{\infty}|p(\xi-j)|^{-2}=0
$$

uniformly on $\Lambda$.
Moving forward, formally define functions $P, Q$ on set the $e^{\Lambda D}$ such that

$$
P\left(e^{\xi D}\right)=\prod_{j=0}^{\infty}|p(\xi+j)| \text { and } Q\left(e^{\xi D}\right)=\prod_{j=1}^{\infty}|p(\xi-j)|^{-1}
$$

Since $X=\overline{e^{\Lambda D}}$ has positive measure in $H$, we may extend the definition of $P$ and $Q$ to $X$ such that $Q=P=0$ on $X$. Let $Z$ be a measurable transversal for $H / G$ in $G$. Since $X$ has positive $\mu_{H}$-measure, there exists a unique $H$-invariant measure $\mu_{Z}$ [11, Theorem 2.51] on $Z$ such that

$$
\begin{equation*}
\mu_{H}(X)=\int_{H} 1_{X}(h) d \mu_{H}(h)=\int_{Z} \int_{G} 1_{X}(z g) d \mu_{G}(g) d \mu_{Z}(z)>0 . \tag{13}
\end{equation*}
$$

Equation (13) gives a disintegration of the Haar measure on $H$ into a family of measures, parameterized by the transversal $Z$ and defined over the cosets of $G$. Additionally, by virtue of (13), there exists at least one coset $z G$ of $G$ such that $z G \cap X$ has positive measure in $z G$. Fixing such a coset, the matrix $h \in z G \cap X$, viewed as an element of $H$ may be factored uniquely such that $h=z g, z \in Z$ and additionally,

$$
P(h)=P(z g)=\prod_{j=0}^{\infty}\left|\left\langle z g e^{j D} c, v\right\rangle\right|=Q(h)=Q(z g)=\prod_{j=1}^{\infty}\left|\left\langle z g e^{-j D} c, v\right\rangle\right|^{-1}=0
$$

for all $g$ in a set $\Omega$ of positive measure in $G$.
Fix $\epsilon \in(0,1)$. Choose a natural number $L$, sufficiently large to ensure that for all $n>L$, and $z g \in z \Omega$,

$$
\left\{\begin{array}{l}
\prod_{j=0}^{n-1}\left|\left\langle z g e^{j D} c, v\right\rangle\right|<\epsilon \\
\prod_{j=1}^{n}\left|\left\langle z g e^{-j D} c, v\right\rangle\right|^{-1}<\epsilon
\end{array}\right.
$$

Since for $z \in Z$, the map $\sigma: z G \rightarrow z G$ defined by $\sigma\left(z_{1} g_{1}\right)=e^{D} z_{1} g_{1}$ is ergodic, and by the Ergodic theorem, there exists $n>L$ such that $\mu_{z G}\left(\sigma^{n}(z \Omega) \cap z \Omega\right)>0$. For $g \in \sigma^{n}(z \Omega) \cap z \Omega$,
we may write $g=z \omega$ and

$$
\prod_{j=1}^{n}\left|\left\langle z \omega e^{-j D} c, v\right\rangle\right|^{-1}=\prod_{j=1}^{n}\left|\left\langle g e^{-j D} c, v\right\rangle\right|^{-1}<\epsilon
$$

for $n>L$. Similarly, since $\sigma^{-n}(g) \in z \Omega$, the following inequality is satisfied:

$$
\prod_{j=0}^{n-1}\left|\left\langle\sigma^{-n}(g) e^{j D} c, v\right\rangle\right|=\prod_{j=0}^{n-1}\left|\left\langle g e^{(j-n) D} c, v\right\rangle\right|<\epsilon
$$

The change of variables $\ell=j-n$ yields

$$
\prod_{\ell=-n}^{-1}\left|\left\langle g e^{\ell D} c, v\right\rangle\right|=\prod_{\ell=1}^{n}\left|\left\langle g e^{-\ell D} c, v\right\rangle\right|<\epsilon
$$

In summary, we have found $n>L$ such that

$$
\left\{\begin{array}{l}
\prod_{j=1}^{n}\left|\left\langle g e^{-j D} c, v\right\rangle\right|^{-1}<\epsilon \\
\prod_{j=1}^{n}\left|\left\langle g e^{-j D} c, v\right\rangle\right|<\epsilon
\end{array}\right.
$$

which is a contradiction, since $0<\epsilon<1$.

## Remark 8.

(a) We say that two configurations $\Lambda_{1}, \Lambda_{2}$ of $n$ distinct points in $\mathbb{R}^{2}$ are equivalent if and only if there exists a symplectic matrix $A \in S L(2, \mathbb{R})$ such that $\Lambda_{2}=A \Lambda_{1}$. This relation defines an equivalence relation on the set of all n-point configurations. As a result, we can partition the collection of all sets of $n$ distinct points in $\mathbb{R}^{2}$ under this equivalence relation. Therefore, when we choose an n-point configuration, we may assume without any loss of generality that it is a representative for a corresponding equivalence class. Furthermore, when $n \geq 3$ we may assume that the chosen (representative) $n$-point configuration includes the set $\{(0,0),(0,1),(a, 0)\}$ where $a \neq 0$, [10, 12, 16].
(b) We now consider the case where $n=4$. By Theorem 6, the assumption $g \in \mathcal{S}(\mathbb{R})$ can be removed from part (2) of Proposition 1. Consequently, to settle the conjecture for sets of 4 points, one needs to consider configurations that are equivalent to $\Lambda_{g}=$ $\{(0,0),(0,1),(a, 0),(\alpha, \beta)\}$ that are neither $(1,3)$-configuration, $(2,2)$-configuration, nor $\left(1, \mathbb{Z}^{2}\right)$-configurations. Let $0 \neq g \in L^{2}(\mathbb{R})$ and suppose that $\Lambda_{g}$ is the collection of all 4-point configurations for which the HRT conjecture fails for this function $g$. The ultimate goal is to prove that such collection is empty. Thus far, we can only say that it does not contain any $(1,3)$-configuration, $(2,2)$-configuration, nor a $\left(1, \mathbb{Z}^{2}\right)-$ configuration. A question that has not received much attention is to understand the
structures of $\Lambda_{g}$. For example, is it true that if $g_{1} \neq g_{2} \in L^{2}(\mathbb{R})$, then $\Lambda_{g_{1}}$ and $\Lambda_{g_{2}}$ are disjoint? In particular, suppose that $\mathcal{G}\left(g, \Lambda_{g}\right)$ is linearly dependent, and let $\Lambda^{\prime}$ be equivalent to $\Lambda_{g}$, then is it the case that $\mathcal{G}\left(g, \Lambda^{\prime}\right)$ is also linearly dependent?

## Acknowledgment

K.A.O. was partially supported by a grant from the National Science Foundation grant DMS 1814253.

## References

1. J. Bak and D. J. Newman, Complex analysis, third ed., Undergraduate Texts in Mathematics, Springer, New York, 2010. MR 2675489
2. J. J. Benedetto and A. Bourouihiya, Linear independence of finite Gabor systems determined by behavior at infinity, J. Geom. Anal. 25 (2015), no. 1, 226-254.
3. M. Bownik and D. Speegle, Linear independence of time-frequency translates of functions with faster than exponential decay, Bull. Lond. Math. Soc. 45 (2013), no. 3, 554-566.
4. C. Demeter, Linear independence of time frequency translates for special configurations, Math. Res. Lett. 17 (2010), no. 4, 761-779.
5. C. Demeter and A. Zaharescu, Proof of the HRT conjecture for $(2,2)$ configurations, J. Math. Anal. Appl. 388 (2012), no. 1, 151-159.
6. M. Einsiedler and T. Ward, Ergodic theory with a view towards number theory, Graduate Texts in Mathematics, vol. 259, Springer-Verlag London, Ltd., London, 2011. MR 2723325
7. H. G. Feichtinger and T. Strohmer (eds.), Gabor analysis: theory and application, Birkhaüser, Boston, MA, 1998.
8. H. G. Feichtinger and T. Strohmer (eds.), Advances in Gabor analysis, Birkhaüser, Boston, MA, 2003.
9. J. Feldman and F. P. Greenleaf, Existence of Borel transversals in groups, Pacific Journal of Mathematics 25 (1968), no. 3, $455-461$.
10. G. B. Folland, Harmonic analysis in phase space, Annals of Mathematics Studies, vol. 122, Princeton University Press, Princeton, NJ, 1989.
11. _ A course in abstract harmonic analysis, vol. 29, CRC press, 2016.
12. K. Gröchenig, Foundations of Time-Frequency Analysis, Applied and Numerical Harmonic Analysis, Springer-Birkhäuser, New York, 2001.
13. B. C. Hall, Quantum theory for mathematicians, Graduate Texts in Mathematics, vol. 267, Springer, New York, 2013. MR 3112817
14. C. Heil, History and evolution of the density theorem for Gabor frames, Journal of Fourier Analysis and Applications 13 (2007), no. 2, 113-166.
15. _ A basis theory primer: expanded edition, Springer Science \& Business Media, 2010.
16. C. Heil, J. Ramanathan, and P. Topiwala, Linear independence of time-frequency translates, Proc. Amer. Math. Soc. 124 (1996), no. 9, 2787-2795.
17. R. Howe, On the role of the Heisenberg group in harmonic analysis, Bulletin (New Series) of the American Mathematical Society 3 (1980), no. 2, $821-843$.
18. P. A. Linnell, von Neumann algebras and linear independence of translates, Proc. Amer. Math. Soc. 127 (1999), no. 11, 3269-3277. MR 1637388
19. P. A. Linnell, M. J. Puls, A. Roman, et al., Linear dependency of translations and square-integrable representations, Banach Journal of Mathematical Analysis 11 (2017), no. 4, 945-962.
20. W. Liu, Short proof of the HRT conjecture for almost every $(1,3)$ configuration, J. Fourier Anal. Appl. 4 (2019), no. 25, 1350-1360.
21. K. A. Okoudjou, Extension and restriction principles for the HRT conjecture, J. Fourier Anal. Appl. 25 (2019), no. 4, 1874-1901.
22. V. Oussa, Frames arising from irreducible solvable actions I, Journal of Functional Analysis 274 (2018), no. 4, 1202-1254.
23. __ Compactly supported bounded frames on Lie groups, Journal of Functional Analysis 277 (2019), no. 6, 1718-1762.
24. P. Walters, An introduction to ergodic theory, Graduate Texts in Mathematics, vol. 79, Springer-Verlag, New York-Berlin, 1982. MR 648108

Department of Mathematics, Tufts University, Medford MA 02131, USA
Email address: Kasso.Okoudjou@tufts.edu
Department of Mathematics, Bridgewater State University, Bridgewater, MA 0235, USA
Email address: vignon.oussa@bridgew.edu

