

# Linear Independence of Time-Frequency Translates in $\mathbb{R}^d$

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Abstract We establish the linear independence of time-frequency translates for functions f on  $\mathbb{R}^d$  having one-sided decay  $\lim_{x \in H, |x| \to \infty} |f(x)| e^{c|x| \log |x|} = 0$  for all c > 0, which do not vanish on an affine half-space  $H \subset \mathbb{R}^d$ .

Keywords Gabor system  $\cdot$  Exponential decay  $\cdot$  HRT conjecture  $\cdot$  Linear independence  $\cdot$  Time-frequency translates

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## **1** Introduction

The Heil–Ramanathan–Topiwala (HRT) conjecture [12] states that time-frequency translates of a non-zero square integrable function f on  $\mathbb{R}^d$  are linearly independent.<sup>1</sup> There have been a few partial results on this conjecture, mostly focusing on finding conditions on  $\Lambda \subset \mathbb{R}^d \times \mathbb{R}^d$  which guarantee that time-frequency translates

$$\mathcal{G}(f,\Lambda) := \left\{ M_a T_b f = e^{2\pi i a \cdot} f(\cdot - b) : (a,b) \in \Lambda \right\}$$

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 $<sup>^1</sup>$  The original HRT conjecture was only for  $\mathbb{R}$ , but the question is also open for higher dimensions.

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along  $\Lambda$  are linearly independent [3,5–7,16]. Other interesting results related to the HRT conjecture can be found in [1,10,11]. In [4], the authors found a one-sided decay condition that guarantees that arbitrary time-frequency shifts are linearly independent.

The goal of this paper is to generalize the main result of the authors [4] to higher dimensions. We point out that the generalization to higher dimensions of linear independence does not always follow expectations. Using the Fourier transform it is easy to see that in  $\mathbb{R}$ , translates of  $L^p$  functions are linearly independent for  $1 \le p < \infty$ . However, the situation in  $\mathbb{R}^d$  is quite different, as all translates of  $L^p$  functions in  $\mathbb{R}^d$  are linearly independent if and only if  $p \le \frac{2d}{d-1}$ , by the results of Edgar and Rosenblatt [8,19].

The main theorem of this paper can be formulated as follows.

**Theorem 1.1** Let H be an affine half-space in  $\mathbb{R}^d$ , i.e.,  $H = \{x \in \mathbb{R}^d : \langle x, v \rangle > a\}$ for some  $v \in \mathbb{R}^d \setminus \{0\}$  and  $a \in \mathbb{R}$ . Let  $f : \mathbb{R}^d \to \mathbb{C}$  be a Lebesgue measurable function which does not vanish almost everywhere on H. Assume that for all c > 0,

$$\lim_{x \in H, \ |x| \to \infty} |f(x)| e^{c|x| \log |x|} = 0.$$
(1.1)

Then, the set  $\mathcal{G}(f, \mathbb{R}^{2d})$  of time-frequency translates of f is linearly independent. That is,  $\mathcal{G}(f, \Lambda)$  is linearly independent for any  $\Lambda \subset \mathbb{R}^{2d}$ .

Note that, unlike the one-dimensional case, we must make the additional assumption that a function f does not vanish on a half-space. This is because in one dimension, functions which vanish on a tail trivially have linearly independent time-frequency shifts. However, if there is a function  $f \in L^2(\mathbb{R})$  with linearly dependent timefrequency shifts, then  $\mathbf{1}_{[-1,0]} \otimes f \in L^2(\mathbb{R}^2)$  vanishes on a half-space  $\{(x_1, x_2) \in \mathbb{R}^2 : x_1 > 0\}$  and has linearly dependent time-frequency shifts. However, it is possible to remove the assumption that f does not vanish on a half-space in Theorem 1.1 provided we weaken the corresponding conclusion. This is shown in Theorem 3.5.

There are several new ingredients employed in extending the one-dimensional result [4, Theorem 1.1] to the higher-dimensional Theorem 1.1. First, we prove a generalization of the Montgomery–Vaughan inequality [17] to higher dimensions, Theorem 2.1, using the theory of Beurling–Selberg extremal functions for Euclidean balls developed by Holt and Vaaler [13]. We also show a higher-dimensional analogue of the Turán–Nazarov inequality, Theorem 2.2, from the corresponding one-dimensional result [18]. Using this we extend the lower bound estimate on products of trigonometric polynomials from [4] to higher dimensions. Then, we establish the key sufficient condition for the linear independence of time-frequency translates, Theorem 3.1, using the concept of an extended half-space, which induces a total order on  $\mathbb{R}^d$ . We also introduce the notion of directional quasi-norm that enables us to prove Theorem 3.3, which provides a sharper version of Theorem 1.1.

#### **2** Useful Facts

In this section, we recall generalizations of the Montgomery–Vaughan inequality and the Turán–Nazarov inequality to  $\mathbb{R}^d$ . We will also provide proofs of the exact inequal-

ities that we need in our development. In Theorem 2.1 we assume that the dimension d is fixed; all constants are allowed to (implicitly) depend on d.

**Theorem 2.1** (Holt, Vaaler) Fix  $d \in \mathbb{N}$ . For every  $\delta > 0$ , there exists R > 0 such that whenever a trigonometric polynomial

$$u(x) = \sum_{j=1}^{m} c_j e^{2\pi i \langle a_j, x \rangle}, \quad c_j \in \mathbb{C}, \ a_j \in \mathbb{R}^d,$$
(2.1)

satisfies  $\min\{|a_j - a_k| : j \neq k\} \ge \delta$ , we have

$$\frac{|\mathbf{B}_{R}(y)|}{2} \sum_{j=1}^{m} |c_{j}|^{2} \le \int_{\mathbf{B}_{R}(y)} |u(x)|^{2} dx \quad \text{for all } y \in \mathbb{R}^{d}.$$
 (2.2)

*Here*,  $|\mathbf{B}_R(y)|$  *is d-dimensional Lebesgue measure of the Euclidean ball*  $\mathbf{B}_R(y)$  *of radius* R *centered at* y.

*Proof* This is an immediate corollary to [13, Theorem 4]. We briefly include the relevant facts for completeness; all references are to [13] and notation is from [13, Theorem 1]. Let  $\xi$ ,  $\delta$ , and  $\nu$  be real numbers with  $\delta > 0$  and  $\nu > -1$ . Define  $u_{\nu}(\xi, \delta)$  to be the infimum of

$$\frac{1}{2}\int_{-\infty}^{\infty} \left(T(x) - S(x)\right) |x|^{2\nu+1} dx,$$

where the infimum is over all pairs of entire functions *S* and *T* of exponential type at most  $2\pi\delta$  such that

$$S(x) \le \operatorname{sgn}(x - \xi) \le T(x)$$
 for all  $x \in \mathbb{R}$ .

By the estimate [13, p. 204], there is a constant A, depending only on  $\nu$ , such that

$$u_{\nu}(\xi,\delta) \le \delta^{-1}\xi^{2\nu+1} \left(1 + \frac{A}{\xi^2\delta^2}\right) \quad \text{whenever } \xi\delta \ge 1.$$
 (2.3)

By [13, Theorem 4] applied for v = (d - 2)/2 we have

$$\omega_{d-1}\left((2\nu+2)^{-1}R^{2\nu+2} - u_{\nu}(R,\delta)\right)\sum_{j=1}^{m}|c_{j}|^{2} \le \int_{\mathbf{B}_{R}(y)}|u(x)|^{2}\,dx.$$
 (2.4)

Here,  $\omega_{d-1}$  is the surface area of the unit sphere  $S^{d-1} \subset \mathbb{R}^d$ . We note that [13, Theorem 4] requires balls to be centered at 0. However, in the special case when  $d = 2\nu + 2$ , it holds more generally for every ball since the exponent in [13, (1.28)] vanishes and translations correspond to unimodular modifications of coefficients  $c_j$ , j = 1, ..., m.

Choose R > 0 such that  $\delta R \ge 1$  and

$$\frac{1}{\delta R} \left( 1 + \frac{A}{R^2 \delta^2} \right) < \frac{1}{2d}$$

It follows by (2.3) that

$$(2\nu+2)^{-1}R^{2\nu+2} - u_{\nu}(R,\delta) \ge (2\nu+2)^{-1}R^{2\nu+2} - \delta^{-1}R^{2\nu+1}\left(1 + \frac{A}{R^2\delta^2}\right)$$
$$\ge R^{2\nu+1}\left(\frac{R}{d} - \frac{R}{2d}\right) = \frac{R^d}{2d}.$$

Combining this with (2.4) and the fact that  $|\mathbf{B}_R(y)| = R^d \omega_{d-1}/d$  yields (2.2).  $\Box$ 

We will also need a higher-dimensional analogue of the Turán–Nazarov inequality [18]. A similar result for  $\mathbb{Z}^d$ -periodic trigonometric polynomials u, which corresponds to the case when  $a_j \in \mathbb{Z}^d$  in (2.1), was considered by Fontes-Merz [9]. Theorem 2.2 also appears in [2, Lemma 12], but without a proof that we provide below.

**Theorem 2.2** (Higher-dimensional Turán–Nazarov inequality) Let u be a trigonometric polynomial of order m as in (2.1). Let E be any measurable subset of positive measure of a ball  $\mathbf{B}_R(y) \subset \mathbb{R}^d$ , R > 0,  $y \in \mathbb{R}^d$ . There exists an absolute and dimensionless constant A such that

$$\sup_{x \in \mathbf{B}_{R}(y)} |u(x)| \le \left( d2^{d} A \frac{|\mathbf{B}_{R}(y)|}{|E|} \right)^{m-1} \sup_{x \in E} |u(x)|.$$
(2.5)

*Proof* Recall that one-dimensional Turán–Nazarov inequality [18] guarantees the existence of an absolute constant *A* such that for any univariate trigonometric polynomial

$$\tilde{u}(r) = \sum_{j=1}^{m} \tilde{c}_j e^{2\pi i \tilde{a}_j r}, \qquad \tilde{c}_j \in \mathbb{C}, \ \tilde{a}_j \in \mathbb{R},$$
(2.6)

and any measurable subset  $\tilde{E}$  of positive measure of an interval  $I \subset \mathbb{R}$  we have

$$\sup_{r\in I} |\tilde{u}(r)| \le \left(\frac{A|I|}{|\tilde{E}|}\right)^{m-1} \sup_{r\in \tilde{E}} |\tilde{u}(r)|.$$
(2.7)

Let  $z_0 \in \mathbf{B}_R(y)$  be a point that achieves the maximum of |u|, i.e.,

$$|u(z_0)| = \sup_{x \in \mathbf{B}_R(y)} |u(x)|.$$

For any direction  $\omega \in S^{d-1} = \{x \in \mathbb{R}^d : |x| = 1\}$  define a ray section of *E* by

$$E_{\omega} = \left\{ r \in [0, \infty) : z_0 + r\omega \in E \right\}.$$

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Let  $\sigma$  be (d-1)-dimensional Lebesgue measure on  $S^{d-1}$ . By the spherical integration formula we have

$$\begin{aligned} |E| &= \int_{\mathbb{R}^d} \mathbf{1}_E(x) dx = \int_{S^{d-1}} \int_0^\infty \mathbf{1}_E(z_0 + r\omega) r^{d-1} dr d\sigma(\omega) \\ &= \int_{S^{d-1}} \int_0^{2R} \mathbf{1}_{E_\omega}(r) r^{d-1} dr d\sigma(\omega) \\ &\leq (2R)^{d-1} \int_{S^{d-1}} |E_\omega| d\sigma(\omega) \leq (2R)^{d-1} \sigma(S^{d-1}) \mathop{\mathrm{ess\,sup}}_{\omega \in S^{d-1}} |E_\omega|. \end{aligned}$$

Since  $|\mathbf{B}_R(y)| = R^d \sigma(S^{d-1})/d$ , there exists  $\omega_0 \in S^{d-1}$  such that  $E_{\omega_0}$  is Lebesgue measurable and

$$\frac{|E_{\omega_0}|}{2R} \ge \frac{1}{d2^d} \frac{|E|}{|\mathbf{B}_R(y)|}.$$
(2.8)

Define a univariate trigonometric polynomial  $\tilde{u}$  by

$$\tilde{u}(r) = u(z_0 + r\omega_0) = \sum_{j=1}^m \tilde{c}_j e^{2\pi i \tilde{a}_j r}, \quad \text{where } \tilde{c}_j = c_j e^{2\pi i \langle a_j, z_0 \rangle}, \quad \tilde{a}_j = \langle a_j, \omega_0 \rangle.$$

Applying (2.7) for  $\tilde{u}$  and  $\tilde{E} = E_{\omega_0} \subset [0, 2R]$ , by (2.8) we have

$$\sup_{x \in \mathbf{B}_{R}(y)} |u(x)| = |u(z_{0})| \leq \sup_{r \in [0,2R]} |\tilde{u}(r)| \leq \left(A \frac{2R}{|E_{\omega_{0}}|}\right)^{m-1} \sup_{r \in E_{\omega_{0}}} |\tilde{u}(r)|$$
$$\leq \left(d2^{d}A \frac{|\mathbf{B}_{R}(y)|}{|E|}\right)^{m-1} \sup_{x \in E} |u(x)|.$$

This proves (2.5).

As a consequence of Theorems 2.1 and 2.2 we obtain the following generalization of [3, Proposition 2.2].

**Proposition 2.3** Let u be a non-zero trigonometric polynomial as in (2.1). Let R > 0. Then there exists a constant C > 0, depending only on u and R, such that

$$\sup_{x \in \mathbf{B}_R(y)} |u(x)| \ge C \quad \text{for all } y \in \mathbb{R}^d.$$
(2.9)

*Proof* Let  $\delta = \min\{|a_j - a_k| : j \neq k\} > 0$ . Let  $R_0 > 0$  be the corresponding radius as in Theorem 2.1. Then,

$$\frac{1}{2}\sum_{j=1}^{m}|c_{j}|^{2} \leq \frac{1}{|\mathbf{B}_{R_{0}}(y)|}\int_{\mathbf{B}_{R_{0}}(y)}|u(x)|^{2}dx \leq \sup_{x\in\mathbf{B}_{R_{0}}(y)}|u(x)|^{2}.$$
(2.10)

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This shows (2.9) when  $R \ge R_0$ . If  $R < R_0$ , then by Theorem 2.2, there exists a constant c > 0 such that

$$\sup_{x \in \mathbf{B}_R(y)} |u(x)| \ge c \sup_{x \in \mathbf{B}_{R_0}(y)} |u(x)|.$$

By (2.10) this again shows (2.9).

### 3 Proof of the Main Theorem

We start by introducing a technical sufficient condition (3.1) for the linear independence of time-frequency translates of a measurable function, which generalizes the one-dimensional condition in [4]. This is a main ingredient in the proof of our main result, Theorem 1.1. Define the space of all Lebesgue measurable functions on the real line by

 $\mathcal{M} = \{ f : \mathbb{R}^d \to \mathbb{C} \text{ is Lebesgue measurable} \}.$ 

As is customary, we shall identify functions in  $\mathcal{M}$  which are equal almost everywhere.

**Definition 3.1** Given an orthonormal basis  $\{v_j\}_{j=1}^d$  of  $\mathbb{R}^d$  define an *extended half-space* by

$$H^{\{v_j\}} = \bigcup_{j=1}^d \{x \in \mathbb{R}^d : \langle x, v_j \rangle > 0 \text{ and } \langle x, v_i \rangle = 0 \text{ for all } i = 1, \dots, j-1\}.$$

Note that the extended half-space  $H^{\{v_j\}}$  essentially coincides with the open half-space

$$H = \{x \in \mathbb{R}^d : \langle x, v_1 \rangle > 0\} \subset H^{\{v_j\}} \subset \overline{H} = \{x \in \mathbb{R}^d : \langle x, v_1 \rangle \ge 0\}.$$

Indeed,  $H^{\{v_j\}} \setminus H \subset \{x \in \mathbb{R}^d : \langle x, v_1 \rangle = 0\}$  has measure zero.

**Theorem 3.1** Let  $H^{\{v_j\}}$  be an extended half-space, and  $H = \{x \in \mathbb{R}^d : \langle x, v_1 \rangle > a\}$ ,  $a \in \mathbb{R}$ , be an affine half-space, and  $f \in \mathcal{M}$ . Suppose that for any non-zero trigonometric polynomial u, any finite subset  $B = \{b_1, \ldots, b_n\} \subset H^{\{v_j\}}$ , and any M > 0, the set

$$E = E_{u,M,B} = \left\{ x \in H : |u(x)f(x)| > M \sum_{i=1}^{n} |f(x+b_i)| \right\}$$
(3.1)

has positive measure. Then,  $\mathcal{G}(f, \mathbb{R}^{2d})$  is linearly independent.

*Proof* Suppose for the sake of contradiction that there exist  $b_1, \ldots, b_N \in \mathbb{R}^d$  and trigonometric polynomials  $u_1, \ldots, u_N$  such that

$$\sum_{i=1}^{N} u_i(x) f(x - b_i) = 0 \quad \text{for a.e. } x \in \mathbb{R}^d.$$

The extended half-space  $H^{\{v_j\}}$  induces a total order  $\prec$  on  $\mathbb{R}^d$  given by

$$x \prec y \iff y - x \in H^{\{v_j\}}$$

This is a consequence of the observation that two extended half-spaces  $H^{\{v_j\}}$  and  $-H^{\{v_j\}} = H^{\{-v_j\}}$  form a partition of  $\mathbb{R}^d \setminus \{0\}$ . Hence, without loss of generality we can assume that

$$b_1 \prec \cdots \prec b_N. \tag{3.2}$$

Moreover, we can also assume that  $||u_i||_{\infty} \leq 1$  for all i = 1, ..., N.

We shall prove that our hypothesis (3.1) implies that there exist sets of positive measure  $Q_1, \ldots, Q_N \subset H$  such that the matrix

$$M_{N} = \begin{pmatrix} u_{1}(x_{1})f(x_{1}-b_{1}) & u_{2}(x_{1})f(x_{1}-b_{2}) & \cdots & u_{N}(x_{1})f(x_{1}-b_{N}) \\ u_{1}(x_{2})f(x_{2}-b_{1}) & u_{2}(x_{2})f(x_{2}-b_{2}) & \cdots & u_{N}(x_{2})f(x_{2}-b_{N}) \\ \vdots & & \vdots \\ u_{1}(x_{N})f(x_{N}-b_{1}) & u_{2}(x_{N})f(x_{N}-b_{2}) & \cdots & u_{N}(x_{N})f(x_{N}-b_{N}) \end{pmatrix}$$

has non-zero determinant for almost all  $(x_1, ..., x_N) \in Q_1 \times \cdots \times Q_N$ . This contradicts our hypothesis that the sum of the rows of *M* are zero almost everywhere.

For each  $1 \le n \le N$ , and  $(x_1, \ldots, x_n) \in (\mathbb{R}^d)^n$ , we consider the principal  $n \times n$  submatrix of  $M_N$  given by

$$M_n = M_n(x_1, \dots, x_n) = \begin{pmatrix} u_1(x_1) f(x_1 - b_1) \cdots u_n(x_1) f(x_1 - b_n) \\ \vdots & \vdots \\ u_n(x_n) f(x_n - b_1) \cdots u_n(x_n) f(x_n - b_n) \end{pmatrix}$$

We will show by induction the existence of sets of positive measure  $Q_1, \ldots, Q_n \subset \mathbb{R}^d$ and positive constants  $c_1, \ldots, c_n$  and  $\delta_1, \ldots, \delta_n$  such that

$$|f(x - b_j)| \le c_n$$
 for a.e.  $x \in \bigcup_{i=1}^n Q_i, \ j = 1, ..., n,$  (3.3)

$$|\det M_n(x_1,\ldots,x_n)| \ge \delta_n$$
 for a.e.  $(x_1,\ldots,x_n) \in Q_1 \times \cdots \times Q_n$ . (3.4)

The base case n = 1 follows trivially from the presence of strict inequality in (3.1). Suppose that (3.3) and (3.4) hold for some  $1 \le n < N$ . Let  $\Sigma$  be the set of all permutations of  $\{1, \ldots, n+1\}$  such that  $\sigma(n+1) \ne n+1$ . Then, for any  $(x_1, \ldots, x_n, x_{n+1}) \in Q_1 \times \cdots \times Q_n \times \mathbb{R}^d$ ,

$$|\det M_{n+1}(x_1, \dots, x_n, x_{n+1})| \\\geq |u_{n+1}(x_{n+1}) f(x_{n+1} - b_{n+1}) \det M_n(x_1, \dots, x_n)| \\- \left| \sum_{\sigma \in \Sigma} \prod_{k=1}^{n+1} u_{\sigma(k)}(x_k) f(x_k - b_{\sigma(k)}) \right| \\\geq \delta_n |u_{n+1}(x_{n+1}) f(x_{n+1} - b_{n+1})| - n! (c_n)^n \sum_{i=1}^n |f(x_{n+1} - b_i)|.$$
(3.5)

The last estimate is a consequence of breaking the sum over  $\sigma \in \Sigma$  with  $\sigma(n+1) = i$ , where  $1 \le i \le n$ . By our hypothesis (3.1), the set

$$E = \left\{ x_{n+1} \in H : |u_{n+1}(x_{n+1} + b_{n+1})f(x_{n+1})| > M \sum_{i=1}^{n} |f(x_{n+1} + (b_{n+1} - b_i))| \right\},\$$

where  $M = 2n!(c_n)^n/\delta_n$ , has positive measure. This is because  $b_{n+1} - b_i \in H^{\{v_j\}}$  for i = 1, ..., n by (3.2). We momentarily set  $Q_{n+1} = b_{n+1} + E$ . Then, by (3.5) we have that for almost every  $(x_1, ..., x_{n+1}) \in Q_1 \times \cdots \times Q_{n+1}$ ,

$$|\det M_{n+1}(x_1,\ldots,x_{n+1})| \ge \frac{\delta_n}{2} |u_{n+1}(x_{n+1})f(x_{n+1}-b_{n+1})| > 0.$$

Thus, by restricting to a (positive measure) subset of  $Q_{n+1}$  if necessary, we can find two constants  $c_{n+1}$ ,  $\delta_{n+1} > 0$  such that (3.3) and (3.4) hold, as desired. This completes the proof of Theorem 3.1.

In order to establish Theorem 1.1 we will need the following lemma about products of trigonometric polynomials, which is a consequence of the Turán–Nazarov inequality. Lemma 3.2 is a straightforward generalization of the one-dimensional result [4, Lemma 3.5].

**Lemma 3.2** Let u be a non-zero trigonometric polynomial, let  $B = \{b_1, \ldots, b_n\} \subset \mathbb{R}^d$  be a finite set, and let R > 0. Then, there exists a constant  $\eta = \eta(u, n, R) > 0$  such that for any  $y \in \mathbb{R}^d$  and any  $k \ge 2$ , there exists a measurable subset  $E \subset \mathbf{B}_R(y)$  with  $|E| > |\mathbf{B}_R(y)|/2$  such that

$$\sup_{x \in E} \sum_{i(1)=1}^{n} \cdots \sum_{i(k)=1}^{n} \left| \prod_{j=1}^{k} u \left( x + \sum_{l=1}^{j} b_{i(l)} \right) \right|^{-1} \le e^{\eta k \log k}.$$
 (3.6)

*Proof* Recall that  $\mathbf{B}_R(y)$  denotes a ball centered at  $y \in \mathbb{R}^d$  with radius R > 0. For any  $b \in \mathbb{R}^d$  and t > 0 we define

$$E_b(t) = \{ x \in \mathbf{B}_R(y) : |u(x+b)| < t \}.$$
(3.7)

By Proposition 2.3 we have (2.9). Combining this with Theorem 2.2 yields

$$t \ge \sup_{x \in E_b(t)} |u(x+b)| \ge C (d2^d A)^{1-m} \left(\frac{|E_b(t)|}{|\mathbf{B}_R(y)|}\right)^{m-1}$$

Thus,

$$|E_b(t)| \le C' |\mathbf{B}_R(y)| t^{1/(m-1)} \quad \text{for all } t > 0,$$
(3.8)

where the constant C' depends on d, R and u, but not on b or y.

For fixed  $k \in \mathbb{N}$  define the set

$$\Sigma = \bigg\{ \sum_{i=1}^{n} \alpha_i b_i : \sum_{i=1}^{n} \alpha_i \le k, \ \alpha_i \in \mathbb{N}_0 \bigg\}.$$

Since the sequence  $(\alpha_1, \ldots, \alpha_n, k - (\alpha_1 + \cdots + \alpha_n))$  represents a partition of k into n + 1 blocks, we have

$$\#|\Sigma| \le \binom{k+n}{k} \le Ck^n.$$
(3.9)

For any subset  $\sigma = \{\sigma(1), \ldots, \sigma(k)\} \subset \Sigma$  of size k, define the function

$$f_{\sigma}(x) = \prod_{i=1}^{k} \frac{1}{|u(x+\sigma(i))|}.$$

Let t > 0. Suppose that for some  $x \in \mathbf{B}_R(y)$  we have

$$\sum_{i(1)=1}^{n} \cdots \sum_{i(k)=1}^{n} \left| \prod_{j=1}^{k} u \left( x + \sum_{l=1}^{j} b_{i(l)} \right) \right|^{-1} > t.$$
(3.10)

By taking averages, this implies that there exists a subset  $\sigma \subset \Sigma$  of size k such that  $f_{\sigma}(x) > t/n^k$ . Since  $f_{\sigma}$  is a product of k functions, at least one of them must take value greater than  $(t/n^k)^{1/k}$ . That is,

$$x \in \bigcup_{i=1}^{k} E_{\sigma(i)}\left(\frac{n}{t^{1/k}}\right) \subset \bigcup_{b \in \Sigma} E_b\left(\frac{n}{t^{1/k}}\right),$$

where  $E_b(t)$  is given by (3.7). Thus, using (3.8) and (3.9), the Lebesgue measure of the set of points  $x \in \mathbf{B}_R(y)$  satisfying (3.10) is bounded by

$$\left|\bigcup_{b\in\Sigma} E_b\left(\frac{n}{t^{1/k}}\right)\right| \le \#|\Sigma| \max_{b\in\Sigma} \left|E_b\left(\frac{n}{t^{1/k}}\right)\right| \le Ck^n C' |\mathbf{B}_R(y)| \frac{n^{1/(m-1)}}{t^{1/k(m-1)}} \le C'' |\mathbf{B}_R(y)| k^n t^{-\frac{1}{k(m-1)}}.$$

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If we wish that the measure of this set does not exceed  $|\mathbf{B}_R(y)|/2$ , we are led to the inequality

$$t > (2C'')^{(m-1)k} k^{n(m-1)k}$$

Thus, there exists a constant  $\eta > 0$ , which is independent of the choice of  $k \ge 2$ , such that  $t = e^{\eta k \log k}$  satisfies the above bound. Consequently, the set *E* of points  $x \in \mathbf{B}_R(y)$  such that the inequality (3.10) **fails** has measure at least  $|\mathbf{B}_R(y)|/2$ . This completes the proof of Lemma 3.2.

We are now ready to state Theorem 3.3. As we will see, the main result of the paper, Theorem 1.1, follows immediately from it.

**Theorem 3.3** Let  $\{v_j\}_{j=1}^d$  be an orthonormal basis of  $\mathbb{R}^d$ . Define the corresponding directional quasi-norm as a mapping  $N : \mathbb{R}^d \to [0, \infty)$  given for  $x \in \mathbb{R}^d$  by

$$N(x) = \sum_{j=1}^{d} \langle x, v_j \rangle_+, \quad \text{where } y_+ = \max(y, 0). \quad (3.11)$$

Let  $f : \mathbb{R}^d \to \mathbb{C}$  be a Lebesgue measurable function that does not vanish almost everywhere on an affine half-space  $H = \{x \in \mathbb{R}^d : \langle x, v_1 \rangle > a\}$ ,  $a \in \mathbb{R}$ . Assume that f satisfies for all c > 0,

$$\lim_{t \to \infty} e^{ct \log t} \sup_{x \in H, N(x) > t} |f(x)| = 0.$$
(3.12)

Then, the set  $\mathcal{G}(f, \mathbb{R}^{2d})$  of time-frequency translates of f is linearly independent.

In the proof of Theorem 3.3 we will need the following lemma.

**Lemma 3.4** Let  $B = \{b_1, \ldots, b_n\}$  be a finite subset of an extended half-space  $H^{\{v_j\}}$ . Then, there exists  $\delta > 0$  such that for any  $k \in \mathbb{N}$  and any choice of  $i(l) \in \{1, \ldots, n\}$ ,  $l = 1, \ldots, k$  we have

$$N\bigg(\sum_{l=1}^k b_{i(l)}\bigg) > \delta k.$$

*Proof* We shall proceed by induction on the dimension *d*. Lemma 3.4 is trivially true when d = 1. Assume by inductive hypothesis that it is true in the dimension d - 1. Without loss of generality, we can assume that elements of *B* are arranged in increasing order  $\prec$  as in the proof of Theorem 3.1, i.e.,  $b_1 \prec \cdots \prec b_n$ . Let  $s = 1, \ldots, n$  be the largest index such that  $\langle b_s, v_1 \rangle = 0$ . If such *s* does not exist, then we let s = 0. Observe that

$$0 < \langle b_{s+1}, v_1 \rangle \le \dots \le \langle b_n, v_1 \rangle. \tag{3.13}$$

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On the other hand, the elements  $\{b_1, \ldots, b_s\}$  lie in the subspace span $\{v_2, \ldots, v_d\}$ , which we can identify with  $\mathbb{R}^{d-1}$ . Since  $H^{\{v_j\}_{j=2}^d}$  is an extended half-space in  $\mathbb{R}^{d-1}$ , by the inductive hypothesis there exists  $\delta > 0$  such that

$$N\left(\sum_{l=1,\ i(l)\leq s}^{k} b_{i(l)}\right) > \delta(k-k_0), \tag{3.14}$$

where  $k_0$  is the number of l = 1, ..., k such that i(l) > s. Moreover, we have

$$N\left(\sum_{l=1}^{k} b_{i(l)}\right) \ge N\left(\sum_{l=1, i(l) \le s}^{k} b_{i(l)}\right) - k_0 dC,$$
(3.15)

where  $C = \max\{|b_i| : i = 1, ..., n\}$ . Indeed, (3.15) follows easily from

$$\left\langle \sum_{l=1}^{k} b_{i(l)}, v_{j} \right\rangle_{+} \ge \left\langle \sum_{l=1, i(l) \le s}^{k} b_{i(l)}, v_{j} \right\rangle_{+} - k_{0}C \quad \text{for all } j = 1, \dots, d.$$

Combining (3.14) and (3.15) we have

$$N\left(\sum_{l=1}^{k} b_{i(l)}\right) > \delta(k-k_0) - k_0 dC \ge \frac{\delta}{2}k, \quad \text{if } k_0 \le \varepsilon k, \quad (3.16)$$

where  $\varepsilon = \delta/(2\delta + 2dC)$ . However, if  $k_0 > \varepsilon k$ , then by (3.13) we have

$$N\left(\sum_{l=1}^{k} b_{i(l)}\right) \ge \left\langle\sum_{l=1}^{k} b_{i(l)}, v_{1}\right\rangle_{+} \ge \left\langle\sum_{l=1, i(l)>s}^{k} b_{i(l)}, v_{1}\right\rangle$$
$$\ge k_{0}\langle b_{s+1}, v_{1}\rangle \ge \varepsilon \langle b_{s+1}, v_{1}\rangle k.$$
(3.17)

Combining (3.16) with (3.17) completes the proof of Lemma 3.4.

*Proof of Theorem 3.3* Let  $H = \{x \in \mathbb{R}^d : \langle x, v \rangle > a\}$  be an affine half-space, where  $v \in \mathbb{R}^d \setminus \{0\}$  and  $a \in \mathbb{R}$ . Choose any orthonormal basis  $\{v_j\}_{j=1}^d \subset \mathbb{R}^d$  such that  $v_1 = v/|v|$ . By Theorem 3.1 it suffices to show that for any trigonometric polynomial  $u \neq 0$ , any finite subset  $B = \{b_1, \ldots, b_n\} \subset H^{\{v_j\}}$  and any M > 0, the set  $E_{u,M,B}$  given by (3.1) has positive measure.

On the contrary, suppose that for some choice of u, B, and M > 0 we have

$$|f(x)| \le M \sum_{i=1}^{n} \frac{|f(x+b_i)|}{|u(x)|}$$
 for a.e.  $x \in H$ . (3.18)

By recursion, (3.18) implies that

$$|f(x)| \le M^k \sum_{i(1)=1}^n \dots \sum_{i(k)=1}^n \left| f\left(x + \sum_{l=1}^k b_{i(l)}\right) \right| \prod_{j=1}^k \left| u\left(x + \sum_{l=1}^{j-1} b_{i(l)}\right) \right|^{-1}.$$
 (3.19)

By the non-vanishing hypothesis on f, there exists a constant  $\varepsilon > 0$ , such that the set  $\{x \in H : |f(x)| > \varepsilon\}$  has positive measure. By the Lebesgue differentiability theorem applied to that set, there exists a ball  $\mathbf{B}_R(y) \subset H$  and such that

$$|\{x \in \mathbf{B}_R(y) : |f(x)| > \varepsilon\}| > |\mathbf{B}_R(y)|/2.$$
(3.20)

Observe that the quasi-norm N defined by (3.11) satisfies the triangle inequality  $N(x + z) \le N(x) + N(z)$ . Thus, for any  $x \in \mathbf{B}_R(y)$  and  $z \in \mathbb{R}^d$ ,

$$N(x + z) \ge N(z) - N(-x) \ge N(z) - \sqrt{d(R + |y|)}$$

By Lemma 3.4 and the Cauchy–Schwarz inequality

$$N(x) \le \sum_{j=1}^{d} |\langle x, v_j \rangle| \le \sqrt{d} \left( \sum_{j=1}^{d} |\langle x, v_j \rangle|^2 \right)^{1/2} = \sqrt{d} |x|,$$
(3.21)

there exists  $\delta > 0$  such that

$$N\left(x+\sum_{l=1}^{k}b_{i(l)}\right) \ge \delta k - \sqrt{d}(R+|y|) \ge \delta k/2 \quad \text{for } k \ge k_0 := 2\sqrt{d}(R+|y|)/\delta.$$

Thus, for any  $x \in \mathbf{B}_R(y)$  and  $k \ge k_0$ ,

$$\left| f\left( x + \sum_{l=1}^{k} b_{i(l)} \right) \right| \le \sup_{z \in H, \ N(z) > \delta k/2} |f(z)|.$$
(3.22)

Here, we used the following fact:  $x \in H$  and  $b \in H^{\{v_j\}} \implies x + b \in H$ .

Combining (3.19) and (3.22) with Lemma 3.2 yields a subset  $E_k \subset \mathbf{B}_R(y)$  with  $|E_k| > |\mathbf{B}_R(y)|/2$  such that

$$|f(x)| \le M^k e^{\eta k \log k} \sup_{z \in H, \ N(z) > \delta k/2} |f(z)| \quad \text{for } x \in E_k.$$
(3.23)

By (3.20) the set  $E_k$  must non-trivially intersect with the set  $\{x \in \mathbf{B}_R(y) : |f(x)| > \varepsilon\}$ . Hence, we conclude that

$$\sup_{z \in H, \ N(z) > \delta k/2} |f(z)| \ge \varepsilon M^{-k} e^{-\eta k \log k} \quad \text{for } k \ge k_0.$$

This contradicts our decay hypothesis (3.12) and completes the proof of Theorem 3.3.

As an immediate consequence of Theorem 3.3 we can deduce Theorem 1.1.

*Proof of Theorem 1.1* Suppose that a function f satisfies the decay condition (1.1). That is, for all c > 0,

$$\lim_{t \to \infty} e^{ct \log t} \sup_{x \in H, |x| > t} |f(x)| = 0.$$

Combining this with (3.21) implies the weaker decay condition (3.12). Consequently, Theorem 3.3 yields the desired conclusion.

We end by presenting a decay condition that is a more direct generalization of the main theorem in [4]. Indeed, the condition (3.24) is automatically satisfied on the real line since any finite subset  $B \subset \mathbb{R}$  can be arranged in increasing order. Then, depending on the sign of  $v \in \mathbb{R}$ , the decay condition (3.25) corresponds to one-sided limit as  $x \to \infty$  or  $x \to -\infty$ . Thus, Theorem 3.5 implies the main result in [4].

**Theorem 3.5** Let  $B = \{b_1, b_2, ..., b_n\}$  be a finite subset of  $\mathbb{R}^d$ . Suppose there exists a vector  $v \in \mathbb{R}^d$  and  $1 \le j_0 \le n$  such that

$$\langle v, b_{j_0} \rangle < \langle v, b_j \rangle$$
 for all  $1 \le j \le n, \ j \ne j_0.$  (3.24)

Let  $f : \mathbb{R}^d \to \mathbb{C}$  be a non-zero Lebesgue measurable function satisfying the directional decay condition

$$\lim_{x \in \mathbb{R}^d, \ \langle x, v \rangle \to \infty} |f(x)| e^{c \langle x, v \rangle \log(\langle x, v \rangle)} = 0 \quad \text{for all } c > 0.$$
(3.25)

Then, if  $u_1, \ldots, u_n$  are trigonometric polynomials such that

$$\sum_{j=1}^{n} u_j(x) f(x+b_j) = 0 \quad \text{for a.e. } x \in \mathbb{R}^d,$$
(3.26)

then  $u_{j_0} = 0$ . In particular, if  $\langle v, b_i \rangle \neq \langle v, b_j \rangle$  for all  $i \neq j$ , then  $\mathcal{G}(f, \mathbb{R}^d \times (-B))$  is linearly independent.

*Remark 3.1* Note that unlike Theorem 1.1 we do not assume that f does not vanish on a half-space. Moreover, the decay condition (1.1) is weakened by the condition (3.25) that does not impose any decay in directions perpendicular to a vector  $v \in \mathbb{R}^d$ . As a consequence, the conclusion of Theorem 3.5 must also be weakened. Indeed, the following simple example shows that we cannot expect that the remaining polynomials satisfy  $u_1 = \cdots = u_n = 0$ . Define  $f : \mathbb{R}^2 \to \mathbb{R}$  by

$$f(x_1, x_2) = \begin{cases} e^{x_2} & \text{if } x_1 \in [0, 1], \\ 0 & \text{otherwise.} \end{cases}$$

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Then, f satisfies the hypothesis of Theorem 3.5 with  $B = \{(0, 0), (0, 1), (-1, 0)\}$ and v = (1, 0), but  $\mathcal{G}(f, \{0\} \times B)$  is linearly dependent.

*Proof of Theorem 3.5* On the contrary suppose that there exists a solution to (3.26) with a non-zero  $u_{j_0}$ . Then,

$$|f(x)| \le M \sum_{j=1, j \ne j_0}^n \frac{|f(x+b_j-b_{j_0})|}{|u_{j_0}(x)|}$$
 for a.e.  $x \in \mathbb{R}^d$ ,

where  $M = \max(||u_1||_{\infty}, ..., ||u_n||_{\infty})$ . Hence, the same inequality as in (3.18) holds true. Define a directional quasi-norm  $\tilde{N}(x) = \langle x, v \rangle_+$ . By the assumption (3.24), Lemma 3.4 holds for the set  $\tilde{B} = \{b_j - b_{j_0} : 1 \le j \ne j_0 \le n\}$  and quasi-norm  $\tilde{N}$  in place of *B* and *N*, resp. Moreover, the same argument as in the proof of Theorem 3.3 works for the quasi-norm  $\tilde{N}$  in place of the original one given by (3.11). As a result we obtain a contradiction with our hypotheses that  $u_{j_0} \ne 0$ , thus completing the proof of Theorem 3.5.

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