ANNALES

## DE

## L'INSTITUT FOURIER

Thomas W. KÖRNER

Fourier Transforms of Measures and Algebraic Relations on Their Supports
Tome 59, n ${ }^{\circ} 4$ (2009), p. 1291-1319.
[http://aif.cedram.org/item?id=AIF_2009__-59_4_1291_0](http://aif.cedram.org/item?id=AIF_2009__-59_4_1291_0)


#### Abstract

© Association des Annales de l'institut Fourier, 2009, tous droits réservés.

L'accès aux articles de la revue «Annales de l'institut Fourier» (http://aif.cedram.org/), implique l'accord avec les conditions générales d'utilisation (http://aif.cedram.org/legal/). Toute reproduction en tout ou partie cet article sous quelque forme que ce soit pour tout usage autre que l'utilisation à fin strictement personnelle du copiste est constitutive d'une infraction pénale. Toute copie ou impression de ce fichier doit contenir la présente mention de copyright.


## cedram

Article mis en ligne dans le cadre du

# FOURIER TRANSFORMS OF MEASURES AND ALGEBRAIC RELATIONS ON THEIR SUPPORTS 

by Thomas W. KÖRNER


#### Abstract

We investigate the relation between the rate of decrease of a Fourier transform and the possible algebraic relations on its support.

Résumé. - Si la transformée de Fourier d'une mesure décroît rapidement alors le support ne satisfait que très peu des relations algèbriques.


## 1. Non-technical introduction

This paper is fairly technical but deals with natural questions.
We work on the circle $\mathbb{T}=\mathbb{R} / \mathbb{Z}$. A well known result tells us that, if a set $E$ has positive Lebesgue measure, then $E+E$ contains an interval. It follows that there exist $x, y \in E$ and integers $m$ and $n$ satisfying some non-trivial equation

$$
m x+n y=0
$$

In other words, if a set has positive Lebesgue measure, it must be rich in short algebraic relations.

A closely related argument shows that any Borel measure $\mu$ whose Fourier transform $\hat{\mu}(r)$ tends fairly rapidly to zero must have a support which is rich in fairly short algebraic relations. More specifically, if $\hat{\mu}(r)=O\left(|r|^{-\epsilon-q^{-1}}\right)$, then we can find $x_{j} \in \operatorname{supp} \mu$ and integers $m_{j}$ satisfying some non-trivial equation

$$
\sum_{j=1}^{q} m_{j} x_{j}=0
$$

Keywords: Convolution, Fourier series.
Math. classification: 42A16.

In an earlier paper, I used a fairly simple probabilistic argument to construct a Borel measure $\mu$ such that $\hat{\mu}(r)=O\left(|r|^{\epsilon-2^{-1} q^{-1}}\right)$, but there do not exist $x_{j} \in \operatorname{supp} \mu$ and integers $m_{j}$ satisfying some non-trivial equation

$$
\sum_{j=1}^{q} m_{j} x_{j}=0
$$

There is a large gap between the results of the two paragraphs and both seem 'natural'. However, in Theorem 2.4, I show that that, by using more complicated probabilistic arguments, we can construct a Borel measure $\mu$ such that $\hat{\mu}(r)=O\left(|r|^{\epsilon-2^{-1} q^{-1}}\right)$, but there do not exist $x_{j} \in \operatorname{supp} \mu$ and integers $m_{j}$ satisfying some non-trivial equation

$$
\sum_{j=1}^{q+1} m_{j} x_{j}=0
$$

If $\epsilon$ is small, the set

$$
\left\{\sum_{j=1}^{q+1} x_{j}: x_{j} \in \operatorname{supp} \mu\right\}
$$

has positive Lebesgue measure and this suggests that the new result is close to best possible or, at least, that it will be quite hard to improve.

On the other hand, if we deal with sets, I show (in Theorem 2.6) how to construct a closed set $E$ such that the $q$-fold sum $E+E+\cdots+E$ has positive Lebesgue measure but there do not exist $x_{j} \in \operatorname{supp} \mu$ and integers $m_{j}$ satisfying some non-trivial equation

$$
\sum_{j=1}^{2 q-1} m_{j} x_{j}=0
$$

## 2. Technical introduction

As stated earlier, we work on the circle $\mathbb{T}=\mathbb{R} / \mathbb{Z}$. All measures will be Borel measures and $m$ will denote the Lebesgue measure. If $\mu$ is a measure, we write

$$
\mu_{[q]}=\mu * \mu * \cdots * \mu
$$

for the $q$-fold convolution of $\mu$ with itself and

$$
E_{[q]}=E+E+\cdots+E=\left\{\sum_{j=1}^{q} x_{j}: x_{j} \in E\right\}
$$

As usual $\|f\|_{1}=\int_{\mathbb{T}}|f(t)| d t$.
This paper, like its predecessor [4], centres round the following two simple observations.

Lemma 2.1. - Suppose that $\mu$ is a non-zero measure with support $E$.
(i) If $\sum_{r=-\infty}^{\infty}|\hat{\mu}(r)|^{q}$ converges, then there exists a non-trivial interval I such that every $x \in I$ can be written

$$
x=x_{1}+x_{2}+\cdots+x_{q}
$$

with $x_{j} \in E$.
(ii) If $\sum_{r=-\infty}^{\infty}|\hat{\mu}(r)|^{2 q}$ converges, then there exists a set $A$ of strictly positive Lebesgue measure such that every $x \in A$ can be written

$$
x=x_{1}+x_{2}+\cdots+x_{q}
$$

with $x_{j} \in E$.
Proof. -
(i) Observe that

$$
\left|\hat{\mu}_{[q]}(r)\right|=|\hat{\mu}(r)|^{q}
$$

so $d \mu_{[q]}=f d m$ where $f$ has an absolutely convergent Fourier series and so is continuous. The support of $f$ contains a non-trivial interval $I$ and

$$
\operatorname{supp} f=\operatorname{supp} \mu_{[q]} \subseteq\left\{x_{1}+x_{2}+\cdots+x_{q}: x_{j} \in E\right\}
$$

(ii) Observe that $d \mu_{[q]}=f d m$ where $f \in L^{2}(m)$ and argue much as in (i).

As I remarked earlier, part (i) of the next lemma is extremely well known, but, although part (ii) is a simple consequence, I do not know if it has been observed before.

Lemma 2.2. - (i) If $E$ has strictly positive Lebesgue measure, then $E+E$ contains a non-trivial interval.
(ii) If $E$ has strictly positive Lebesgue measure, then we can find a nontrivial interval $I$ such that, whenever $x \in I$, the equation

$$
x_{1}+x_{2}=x
$$

has uncountably many distinct solutions with $x_{1}, x_{2} \in E$.

Proof. -
(i) If $E$ has strictly positive Lebesgue measure then we can find a closed set $E^{*} \subseteq E$ with $E^{*}$ having strictly positive measure. Thus, without loss of generality, we may assume that $E$ is closed. We now know that the indicator function $\mathbb{I}_{E}$ is a nontrivial $L^{2}(m)$ function so $\sum_{j=-\infty}^{\infty}\left|\hat{\mathbb{I}}_{E}(j)\right|^{2}$ converges and we may apply Lemma 2.1 (i).
(ii) Suppose that the result is false. Then each interval $I$ contains a point $y$ such that the equation

$$
x_{1}+x_{2}=y
$$

has only countably many distinct solutions with $x_{1}, x_{2} \in E$. Thus we can find a countable dense sequence $y_{j}$ and associated countable sets $E_{j}$ such that, if

$$
x_{1}+x_{2}=y_{j}
$$

with $x_{1}, x_{2} \in E$ then $x_{1}, x_{2} \in E_{j}$. Now observe that $E \backslash \bigcup_{j=1}^{\infty} E_{j}$ is a set of strictly positive Lebesgue measure disobeying the conclusions of (i) which is impossible

Since every non-trivial interval contains a rational, Lemma 2.1 (i) implies the following result.

Lemma 2.3. - Suppose that $\mu$ is a non-zero measure on $\mathbb{T}$ and $q$ is a positive integer such that we can find an $\alpha>1 / q$ and an $A>0$ with

$$
|\hat{\mu}(r)| \leqslant A|r|^{-\alpha}
$$

for all $r \neq 0$. Then we can find distinct points $x_{1}, x_{2}, \ldots, x_{q} \in \operatorname{supp} \mu$ and $m_{j} \in \mathbb{Z}$, not all zero, such that

$$
\sum_{j=1}^{q} m_{j} x_{j}=0
$$

In this paper we show how to prove the following result in the other direction.

Theorem 2.4. - If $q$ is an integer with $q \geqslant 1$ and $\psi: \mathbb{N} \rightarrow \mathbb{R}$ is a sequence of strictly positive numbers such that $\psi(r) \rightarrow \infty$ as $r \rightarrow \infty$, then there exists a probability measure $\mu$ such that

$$
|\hat{\mu}(r)| \leqslant|r|^{-1 /(2 q)}(\log (1+|r|))^{1 / 2} \psi(|r|)
$$

for all $r \neq 0$, but, given distinct points $x_{1}, x_{2}, \ldots, x_{q+1} \in \operatorname{supp} \mu$, the only solution to the equation

$$
\sum_{j=1}^{q+1} m_{j} x_{j}=0
$$

with $m_{j} \in \mathbb{Z}$ is the trivial solution $m_{1}=m_{2}=\cdots=m_{q+1}=0$.
In [4] we proved a similar result with the equation $\sum_{j=1}^{q+1} m_{j} x_{j}=0$ replaced by $\sum_{j=1}^{q} m_{j} x_{j}=0$. Earlier I explained why the new result might be substantially more difficult to prove than the old. Observe that if, for example, $\psi(r)=(\log (1+|r|))^{1 / 2}$, then, by Lemma 2.1,

$$
\left\{\sum_{j=1}^{q+1} x_{j}: x_{j} \in \operatorname{supp} \mu\right\}
$$

must have strictly positive Lebesgue measure.
The key lemma in our proof is the following.
Lemma 2.5. - Let $q$ be an integer with $q \geqslant 1$ and $\psi: \mathbb{N} \rightarrow \mathbb{R}$ be a sequence of strictly positive numbers such that $\psi(r) \rightarrow \infty$ as $r \rightarrow \infty$.

Suppose $\mathbf{m} \in \mathbb{Z}^{q+1} \backslash\{\mathbf{0}\}$ and $N$ is a positive integer such that

$$
N \geqslant 12\left(q+1+\sum_{j=1}^{N}\left|m_{j}\right|\right)
$$

Then, given closed intervals $I_{j}=\left[\left(n_{j}-\frac{1}{2}\right) / N,\left(n_{j}+\frac{1}{2}\right) / N\right]$, with $n_{j}$ an integer, such that

$$
\left|\frac{n_{j}}{N}-\frac{n_{k}}{N}\right| \geqslant \frac{6}{N} \text { for } 1 \leqslant j<k \leqslant q+1
$$

and $\epsilon>0$, we can find an infinitely differentiable function $f$ with the following properties.
(i) $f(t) \geqslant 0$ for all $t \in \mathbb{T}$.
(ii) $\hat{f}(0)=1$.
(iii) $|\hat{f}(r)| \leqslant|r|^{-1 /(2 q)}(\log (1+|r|))^{1 / 2} \psi(|r|)$ for all $r \neq 0$.
(iv) If $x_{j} \in \operatorname{supp} f \cap I_{j}$ then

$$
\sum_{j=1}^{q+1} m_{j} x_{j} \neq 0
$$

(v) If $x \in \mathbb{T}$ we can find a $y \in \operatorname{supp} f$ with $|x-y|<\epsilon$.

If we deal with sets rather than measures we have the following result which excludes a natural conjecture.

THEOREM 2.6. - If $q$ is an integer with $q \geqslant 1$, then we can find a closed set $E$ with the following properties.
(i) $E_{[q]}$ has strictly positive Lebesgue measure.
(ii) The equation

$$
\sum_{j=1}^{2 q-1} m_{j} x_{j}=0
$$

has no non-trivial solution with $m_{j} \in \mathbb{Z}$ and the $x_{j}$ distinct points of $E$.

The $\mu$ we construct in the proof of Theorem 2.4 also has the property described in the next lemma, which furnishes a complement to Lemma 2.1 (ii).

Lemma 2.7. - If $q$ is an integer with $q \geqslant 1$ and $\psi: \mathbb{N} \rightarrow \mathbb{R}$ is a sequence of positive numbers such that $\psi(r) \rightarrow \infty$ as $r \rightarrow \infty$, then there exists a probability measure $\mu$ such that

$$
|\hat{\mu}(r)| \leqslant|r|^{-1 /(2 q)}(\log (1+|r|))^{1 / 2} \psi(|r|)
$$

for all $r \neq 0$, but the set

$$
\left\{\sum_{j=1}^{q} m_{j} x_{j}: x_{j} \in \operatorname{supp} \mu, m_{j} \in \mathbb{Z}\right\}
$$

has Lebesgue measure zero.
However, the method of [4] can be easily adapted to give a much simpler proof of this result.

Since the proof of Theorem 2.6 is substantially simpler than that of Theorem 2.4 we shall devote the next two sections to its proof. We give the fairly routine proof of Theorem 2.4 from Lemma 2.5 in section 5 and devote the rest of the paper to the proof of Lemma 2.5.

Like many others of my papers, this one owes a great deal to two remarkable papers [2] and [3] of Kaufman.

## 3. Sums and algebraic relations

We shall prove Theorem 2.6 by a Baire category argument. We use the Hausdorff metric $d_{\mathcal{F}}$ defined in the next lemma.

Definition 3.1. - Consider the space $\mathcal{F}$ of non-empty closed subsets of $\mathbb{T}$. We set

$$
d_{\mathcal{F}}(E, F)=\sup _{e \in E} \inf _{f \in F}|e-f|+\sup _{f \in F} \inf _{e \in E}|e-f| .
$$

It is well known that $\left(\mathcal{F}, d_{\mathcal{F}}\right)$ is a complete metric space. (See, for example, Chapter II §21 VII and Chapter III §33 IV of [5].)

We need he following remarks.
Lemma 3.2. - (i) If $E, F, G$ and $H$ are closed then

$$
d_{\mathcal{F}}(E+F, G+H) \leqslant d_{\mathcal{F}}(E, G)+d_{\mathcal{F}}(F, H)
$$

(ii) Suppose $E_{n}, F_{n}, E$ and $F$ are closed sets with

$$
d_{\mathcal{F}}\left(E_{n}, E\right), d_{\mathcal{F}}\left(F_{n}, F\right) \rightarrow 0
$$

Then $d_{\mathcal{F}}\left(E_{n}+F_{n}, E+F\right) \rightarrow 0$ as $n \rightarrow \infty$.
(iii) Suppose $E_{n}$ and $E$ are closed sets with $d_{\mathcal{F}}\left(E_{n}, E\right) \rightarrow 0$. Then

$$
m(E) \geqslant \limsup _{n \rightarrow \infty} m\left(E_{n}\right)
$$

Proof. -
(i) Observe that, if $e \in E, f \in F, g \in G$ and $h \in H$,

$$
|(e+f)-(g+h)| \leqslant|e-g|+|f-h|
$$

so, if $e \in E, f \in F$,

$$
\inf _{e \in E, f \in F}|(e+f)-(g+h)| \leqslant \inf _{e \in E}|e-g|+\inf _{f \in F}|f-h|
$$

whence

$$
\sup _{g \in G, h \in H} \inf _{e \in E, f \in F}|(e+f)-(g+h)| \leqslant \sup _{g \in G} \inf _{e \in E}|e-g|+\sup _{h \in H} \inf _{f \in F}|f-h| .
$$

(ii) This follows directly from (i).
(iii) Given $\epsilon>0$, we can find an $\eta>0$ such that

$$
m(E+(-\eta, \eta))<m(E)+\epsilon
$$

When $n$ is sufficiently large,

$$
E_{n} \subseteq E+(-\eta, \eta)
$$

so

$$
m\left(E_{n}\right)<m(E)+\epsilon .
$$

Thus $\lim \sup _{n \rightarrow \infty} m\left(E_{n}\right) \leqslant m(E)+\epsilon$ for all $\epsilon$ and the result follows.

Definition 3.3. - If $q \geqslant 1$, we define $\mathcal{H}=\mathcal{H}_{q}$ to be the subspace of $\mathcal{F}$ consisting of those closed sets for which $m\left(E_{[q]}\right) \geqslant 1 / 2$ with the inherited metric $d_{\mathcal{H}}=d_{\mathcal{F}} \mid \mathcal{H} \times \mathcal{H}$.

Lemma 3.2 tells us that $\mathcal{H}$ is a closed subspace of $\mathcal{F}$ and so $\left(\mathcal{H}, d_{\mathcal{H}}\right)$ is a complete metric space. Since $(\mathbb{T}, m) \in \mathcal{H}$, the space $\mathcal{H}$ is non-empty.

We can thus deduce Theorem 2.6 from the Baire category version.
Theorem 3.4. - The collection of $E \in \mathcal{H}$ such that the equation

$$
\sum_{j=1}^{2 q-1} m_{j} x_{j}=0
$$

has a non-trivial solution with $m_{j} \in \mathbb{Z}$ and the $x_{j}$ distinct points of $E$ is of first category.

Since $\mathbb{Z}^{2 q-1}$ is countable, Theorem 3.4 follows in turn from the simpler result.

Lemma 3.5. - Let $\mathbf{m} \in \mathbb{Z}^{2 q-1} \backslash\{\mathbf{0}\}$ and $N \geqslant 1$. Let $\mathcal{E}(\mathbf{m}, N)$ be the collection of of $E \in \mathcal{H}$ such that the equation

$$
\sum_{j=1}^{2 q-1} m_{j} x_{j}=0
$$

has a solution with $x_{j} \in E[1 \leqslant j \leqslant 2 q-1]$ and

$$
\left|x_{j}-x_{k}\right| \geqslant N^{-1} \text { for } 1 \leqslant j<k \leqslant 2 q-1
$$

Then $\mathcal{E}(\mathbf{m}, N)$ is closed and has dense complement.
We split the proof of Theorem 3.4 into two parts, the easy Lemma 3.6 and the harder Lemma 3.7.

Lemma 3.6. - Suppose $\mathbf{m} \in \mathbb{Z}^{2 q-1} \backslash\{\mathbf{0}\}$ and $N \geqslant 1$. Then $\mathcal{E}(\mathbf{m}, N)$ is open in $\mathcal{H}$.

Proof. - We show that the complement of $\mathcal{E}(\mathbf{m}, N, \mathbf{I})$ is closed. Suppose that $E_{n} \notin \mathcal{E}(\mathbf{m}, N, \mathbf{I})$ and $d\left(E_{n}, E\right) \rightarrow 0$ as $n \rightarrow \infty$ Then we can find $\mathbf{x}(n) \in E_{n}^{2 q-1}$ such that $\left|x_{j}(n)-x_{i}(n)\right| \geqslant N^{-1}$ for all $i \neq j$ and

$$
\sum_{j=1}^{2 q-1} m_{j} x_{j}(n)=0
$$

The Bolzano-Weierstrass theorem tells us that, by extracting a subsequence if necessary, we may suppose that $x_{j}(n) \rightarrow x_{j}$ as $n \rightarrow \infty$ for each $1 \leqslant j \leqslant$ $2 q-1$. Now $\left|x_{j}-x_{i}\right| \geqslant N^{-1}$ for all $i \neq j$ and

$$
\sum_{j=1}^{2 q-1} m_{j} x_{j}=0
$$

Thus $E_{n} \notin \mathcal{E}(\mathbf{m}, N)$.

Thus our proof of Theorem 3.4 reduces to proving the following result.
Lemma 3.7. - Let $\mathbf{m} \in \mathbb{Z}^{2 q-1} \backslash\{\mathbf{0}\}$, let $\delta>0$ and let $\epsilon>0$. Then, given $E \in \mathcal{H}$ we can find an $F \in \mathcal{H}$ with $d_{\mathcal{H}}(E, F)<\epsilon$ such that the equation

$$
\sum_{j=1}^{2 q-1} m_{j} x_{j}=0
$$

has no solution with $x_{j} \in E[1 \leqslant j \leqslant 2 q-1]$ and

$$
\left|x_{j}-x_{k}\right| \geqslant \delta \text { for } 1 \leqslant j<k \leqslant 2 q-1
$$

## 4. Completion of the proof of the theorem on sums

The main step in the construction for Lemma 3.7 is the following.
Lemma 4.1. - Suppose $E_{1}, E_{2}, \ldots E_{q}$ are closed subsets of $\mathbb{T}$ such that

$$
m\left(E_{1}+E_{2}+\cdots+E_{q}\right) \geqslant 1 / 2
$$

Then, given $\epsilon_{1}>0$ and $\mathbf{m} \in \mathbb{Z}^{2 q-1} \backslash\{\mathbf{0}\}$, we can find $F_{1}, F_{2}, \ldots F_{q}$ closed subsets of $\mathbb{T}$ and $\epsilon_{2}>0$ with the following properties.
(i) $m\left(F_{1}+F_{2}+\cdots+F_{q}\right) \geqslant 1 / 2$.
(ii) The Hausdorff distance $d_{\mathcal{F}}\left(E_{j}, F_{j}\right)<\epsilon_{1}$ for $1 \leqslant j \leqslant q$.
(iii) If $x_{1} \in \bigcup_{s=1}^{q} F_{s} x_{r} \in \bigcup_{s=2}^{q} F_{s}$ for $2 \leqslant r \leqslant 2 q-1,\left|x_{r}-y_{r}\right| \leqslant \epsilon_{2}$ for $1 \leqslant r \leqslant 2 q-1$ and $\left|y_{j}-y_{k}\right| \geqslant \delta$ for $1 \leqslant j<k \leqslant 2 q-1$, we have

$$
\sum_{j=1}^{2 q-1} m_{j} y_{j} \neq 0
$$

Proof. - Choose $0<\gamma<\epsilon_{1} / 4$. We observe that the collection of open sets

$$
\left(E_{1}+(\gamma,-\gamma)\right)+E_{2}+\cdots+E_{q}
$$

form an open cover of the compact set

$$
E_{1}+E_{2}+\cdots+E_{q-1}+E_{q} .
$$

We can thus find a finite collection of points

$$
\mathbf{e}(r) \in E_{2} \times E_{3} \times \cdots \times E_{q} \quad[1 \leqslant r \leqslant N
$$

such that

$$
\bigcup_{r=1}^{N}\left(\sum_{j=2}^{q} e_{j}(r)+\left(E_{1}+(\gamma,-\gamma)\right)\right) \supseteq E_{1}+E_{2}+\cdots+E_{q-1}+E_{q}
$$

We now choose finite subsets $\tilde{F}_{j}$ of $E_{j}[2 \leqslant j \leqslant q]$ such that

$$
\mathbf{e}(r) \in \prod_{j=2}^{q} \tilde{F}_{j}
$$

for $[1 \leqslant r \leqslant N]$ and $d_{\mathcal{F}}\left(E_{j}, \tilde{F}_{j}\right)<\gamma$. Automatically

$$
\left(E_{1}+(\gamma,-\gamma)\right)+\tilde{F}_{2}+\cdots+\tilde{F}_{q} \supseteq E_{1}+E_{2}+\cdots+E_{q-1}+E_{q}
$$

By perturbing each of the points in the $\tilde{F}_{j}$ in turn by an amount less than $\gamma$ we can find disjoint finite sets $F_{j}[2 \leqslant j \leqslant q]$ such that $d_{\mathcal{F}}\left(E_{j}, F_{j}\right)<2 \gamma$,
(iii)' If $y_{r} \in \bigcup_{s=2}^{q} F_{s}$ for $1 \leqslant r \leqslant 2 q-1$ and the $y_{r}$ are distinct, we have

$$
\sum_{j=1}^{2 q-1} m_{j} y_{j} \neq 0
$$

and

$$
\left(E_{1}+(2 \gamma,-2 \gamma)\right)+F_{2}+\cdots+F_{q} \supseteq E_{1}+E_{2}+\cdots+E_{q-1}+E_{q}
$$

A simple argument shows that

$$
m\left(\left(E_{1}+[3 \gamma,-3 \gamma]\right)+F_{2}+\cdots+F_{q}\right)>1 / 2
$$

Since $F_{2}, F_{3}, \ldots, F_{q}$ are finite, it follows that there is a finite set $X$ such that if

$$
y \notin X \text { and } y_{r} \in \bigcup_{s=2}^{q} F_{s} \text { for } 2 \leqslant r \leqslant 2 q-1
$$

with the $y_{r}$ distinct then

$$
m_{1} y+\sum_{j=2}^{2 q-1} m_{j} y_{j} \neq 0
$$

Now set

$$
F_{1}=\left(E_{1}+[3 \gamma,-3 \gamma]\right) \backslash(X+(-\eta, \eta))
$$

with $\eta>0$. Provided we take $\eta$ small enough, we have $d_{\mathcal{F}}\left(E_{1}, F_{1}\right)<\epsilon_{1}$ and

$$
m\left(F_{1}+F_{2}+\cdots+F_{q}\right) \geqslant 1 / 2 .
$$

Further, combining (iii)' with the definition of $X$ we see that
(iii) ${ }^{\prime \prime}$ If $y_{1} \in \bigcup_{s=1}^{q} F_{s} y_{r} \in \bigcup_{s=2}^{q} F_{s}$ for $2 \leqslant r \leqslant 2 q-1$ and $y_{1}, y_{2}, \ldots$, $y_{2 q-1}$ are distinct then

$$
\sum_{j=1}^{2 q-1} m_{j} y_{j} \neq 0
$$

Take $K$ to be the collection of $\mathbf{x} \in \mathbb{T}^{2 q-1}$ such that $x_{1} \in \bigcup_{s=1}^{q} F_{s}$, $x_{r} \in \bigcup_{s=2}^{q} F_{s}$ for $2 \leqslant r \leqslant 2 q-1$ and $\left|x_{j}-x_{k}\right| \geqslant \delta / 2$ for $1 \leqslant j<k \leqslant 2 q-1$. If we set

$$
L=\left\{\mathbf{x} \in \mathbb{T}^{2 q-1}: \sum_{j=1}^{2 q-1} m_{j} x_{j}=0\right\}
$$

then $K$ and $L$ are disjoint compact subsets of $\mathbb{T}^{2 q-1}$. A standard theorem now tells us that there exists an $\epsilon_{2}^{\prime}>0$ such that, if $\mathbf{x} \in K$ and $\left|y_{j}-x_{j}\right| \leqslant \epsilon_{2}^{\prime}$ for $1 \leqslant j \leqslant 2 q-1$, we have

$$
\sum_{j=1}^{2 q-1} m_{j} y_{j} \neq 0
$$

If we take $\epsilon_{2}=\min \left(\epsilon_{2}^{\prime}, \delta\right) / 4$, then condition (iii) holds and the required result follows.

We now show how to prove Lemma 3.7 from Lemma 4.1.
We first observe that, by repeated application of Lemma 4.1, with the various $2 q-1$ tuples obtained by permuting the entries of $\mathbf{m}$ we obtain the following version.

Lemma 4.2. - The result of Lemma 4.1 holds with condition (iii) of the conclusion replaced by the following.
(iii)' Suppose $1 \leqslant p \leqslant 2 q-1$. If $x_{p} \in \bigcup_{s=1}^{q} F_{s}, x_{r} \in \bigcup_{s=2}^{q} F_{s}$ for $r \neq p$, $\left|x_{r}-y_{r}\right| \leqslant \epsilon_{2}$ for $1 \leqslant r \leqslant 2 q-1$, and $\left|y_{j}-y_{k}\right| \geqslant \delta$ for $1 \leqslant j<k \leqslant 2 q-1$, we have

$$
\sum_{j=1}^{2 q-1} m_{j} y_{j} \neq 0
$$

Next observe that, by repeated application of Lemma 4.2 we can obtain the following version.

Lemma 4.3. - Suppose $E_{1}, E_{2}, \ldots E_{q}$ are closed subsets of $\mathbb{T}$ such that

$$
m\left(E_{1}+E_{2}+\cdots+E_{q}\right) \geqslant 1 / 2
$$

Then, given $\epsilon>0, \delta>0$ and $\mathbf{m} \in \mathbb{Z}^{2 q-1} \backslash\{\mathbf{0}\}$, we can find $F_{1}, F_{2}, \ldots F_{q}$ closed subsets of $\mathbb{T}$ and $\eta>0$ with the following properties.
(i) $m\left(F_{1}+F_{2}+\cdots+F_{q}\right) \geqslant 1 / 2$.
(ii) The Hausdorff distance $d_{\mathcal{F}}\left(E_{j}, F_{j}\right)<\epsilon$ for $1 \leqslant j \leqslant q$.
(iii) Suppose $1 \leqslant p_{1} \leqslant 2 q-1$ and $1 \leqslant q_{1} \leqslant q$. If $x_{p_{1}} \in \bigcup_{s=1}^{q} F_{s} x_{r} \in$ $\bigcup_{s \neq q_{1}} F_{s}$ for $r \neq p_{1}$ and $\left|x_{j}-x_{k}\right| \geqslant \delta$ for $1 \leqslant j<k \leqslant 2 q-1$, we have

$$
\sum_{j=1}^{2 q-1} m_{j} x_{j} \neq 0
$$

We now disentangle condition (iii) of Lemma 4.3.
Lemma 4.4. - Condition (iii) of Lemma 4.3 can be rewritten as follows.
(iii)" If $x_{j} \in \bigcup_{s=1}^{q} F_{s}$ for $1 \leqslant j \leqslant 2 q-1$ and $\left|x_{j}-x_{k}\right| \geqslant \delta$ for $1 \leqslant j<$ $k \leqslant 2 q-1$, we have

$$
\sum_{j=1}^{2 q-1} m_{j} x_{j} \neq 0
$$

Proof. - By a simple counting argument there exists a $1 \leqslant q_{1} \leqslant q$ such that $F_{q_{1}}$ contains at most one of the $x_{j}$.

We can now prove Lemma 3.7.
Proof of Lemma 3.7 from Lemma 4.3. - Suppose $E \in \mathcal{H}$. If we set

$$
E_{1}=E_{2}=\cdots=E_{q}=E
$$

then, automatically, $E_{1}, E_{2}, \ldots E_{q}$ are closed subsets of $\mathbb{T}$ such that

$$
m\left(E_{1}+E_{2}+\cdots+E_{q}\right) \geqslant 1 / 2
$$

and so, by Lemma 4.3 (supplemented by the observation of Lemma 4.4), we can find $F_{1}, F_{2}, \ldots F_{q}$ closed subsets of $\mathbb{T}$ with the following properties.
(i) $m\left(F_{1}+F_{2}+\cdots+F_{q}\right) \geqslant 1 / 2$.
(ii) The Hausdorff distance $d_{\mathcal{F}}\left(E_{j}, F_{j}\right)<\epsilon / q$ for $1 \leqslant j \leqslant q$.
(iii) If $x_{j} \in \bigcup_{s=1}^{q} F_{s}$ for $1 \leqslant j \leqslant 2 q-1$ and $\left|x_{j}-x_{k}\right| \geqslant \delta$ for $1 \leqslant j<k \leqslant$ $2 q-1$, we have

$$
\sum_{j=1}^{2 q-1} m_{j} x_{j} \neq 0
$$

If we now set $F=\bigcup_{s=1}^{q} F_{s}$, then simple estimates (not the best possible) give

$$
d_{\mathcal{F}}(E, F) \leqslant \sum_{s=1}^{q} d_{\mathcal{F}}\left(E, F_{s}\right)=\sum_{s=1}^{q} d_{\mathcal{F}}\left(E_{s}, F_{s}\right)<\epsilon
$$

and the remaining conclusions can be read off.

The following observation may be worth making.
Lemma 4.5. - If $E$ is as in Theorem 2.6 then the set

$$
\left\{\sum_{j=1}^{q-1} m_{j} x_{j}: m_{j} \in \mathbb{Z}, x_{j} \in E\right\}
$$

has Lebesgue measure zero.

Proof. - If not, we can find an $\mathbf{m} \in \mathbb{Z}^{q-1}$ such that

$$
F=\left\{\sum_{j=1}^{q-1} m_{j} x_{j}: x_{j} \in E\right\}
$$

has positive Lebesgue measure. But then $F+F$ contains a non-trivial interval which is impossible.

## 5. Proof of the main theorem from the main lemma

Although we have simply demanded that $\psi(r) \rightarrow \infty$ as $r \rightarrow \infty$ in Theorem 2.4, we can demand rather better behaviour.

Lemma 5.1. - If $\tilde{\psi}: \mathbb{N} \rightarrow \mathbb{R}$ is a sequence of strictly positive numbers such that $\tilde{\psi}(r) \rightarrow \infty$ as $r \rightarrow \infty$ and $\delta>0$ we can find an increasing sequence of strictly positive numbers $\psi(r) \rightarrow \infty$ such that
(i) $\min (\tilde{\psi}(r), \delta) \geqslant \psi(r)$ for all $r \in \mathbb{N}$,
(ii) $2 \psi(n) \geqslant \psi(r) \geqslant \psi(n)$ for all $2 n \geqslant r \geqslant n \geqslant 1$.

Proof. - Immediate.
Throughout the rest of this paper $q$ is a fixed integer with $q \geqslant 1$ and $\psi: \mathbb{N} \rightarrow \mathbb{R}$ is a fixed sequence of strictly positive numbers obeying the conditions of Lemma 5.1 with

$$
\delta=\left(\max _{r \geqslant 1}|r|^{-1 /(2 q)}(\log (1+|r|))^{1 / 2}\right)^{-1}
$$

We write

$$
\phi(r)=|r|^{-1 /(2 q)}(\log (1+|r|))^{1 / 2} \psi(|r|)
$$

Observe that $0<\phi(r) \leqslant 1$ for all $r \geqslant 1$ and there is a constant $K \geqslant 1$ such that

$$
K \phi(n) \geqslant \phi(r) \geqslant K^{-1} \phi(n)
$$

for all $2 n \geqslant r \geqslant n \geqslant 1$.
Once again we use a Baire category argument but our metric space is a little more complicated than the Hausdorff metric space $\left(\mathcal{F}, D_{\mathcal{F}}\right)$.

Definition 5.2. - We take $\mathcal{G}$ to be the set of ordered pairs $(E, \mu)$ where $E$ is a non-empty closed set, and $\mu$ is a probability measure such that
(i) $E \supseteq \operatorname{supp} \mu$.
(ii) $|\hat{\mu}(r)| \phi(|r|)^{-1} \rightarrow 0$ as $r \rightarrow \infty$.

If $(E, \mu),(F, \tau) \in \mathcal{G}$ we define

$$
d_{\mathcal{G}}((E, \mu),(F, \tau))=d_{\mathcal{F}}(E, F)+\sup _{r \in \mathbb{Z}}|\hat{\mu}(r)-\hat{\tau}(r)| \phi(|r|)^{-1}
$$

It is easy to check that $\left(\mathcal{G}, d_{\mathcal{G}}\right)$ is a complete metric space. Since $(\mathbb{T}, m) \in$ $\mathcal{G}$, the space is non-empty.

Theorem 2.4 thus follows from its Baire category version.
Theorem 5.3. - The set $\mathcal{E}$ of $(E, \mu)$ such that there exist distinct points

$$
x_{1}, x_{2}, \ldots, x_{q+1} \in E
$$

and integers $m_{j}$, not all zero such that

$$
\sum_{j=1}^{q+1} m_{j} x_{j}=0
$$

is of first category in $\left(\mathcal{G}, d_{\mathcal{G}}\right)$.
It may be worth remarking that we shall use Baire category, not because it gives an apparently more general theorem, but because it makes the book keeping aspects of the proof rather easier. It should also be said that, even if the arguments of this section appear complicated, they are not deep.

In order to attack Theorem 5.3, we introduce some temporary definitions reflecting the conditions of Lemma 2.5. If $\mathbf{m} \in \mathbb{Z}^{q+1} \backslash\{\mathbf{0}\}$ we write

$$
N_{0}(\mathbf{m})=12\left(q+1+\sum_{j=1}^{q+1}\left|m_{j}\right|\right) .
$$

If $N \geqslant 24(q+1)$ we write $\mathcal{J}(N)$ for the collection of ordered $(q+1)$ tuples

$$
\mathbf{I}=\left(I_{1}, I_{2}, \ldots, I_{q+1}\right)
$$

where $I_{j}=\left[\left(n_{j}+\frac{1}{2}\right) / N,\left(n_{j}-\frac{1}{2}\right) / N\right]$, with $n_{j}$ an integer and

$$
\left|\frac{n_{j}}{N}-\frac{n_{k}}{N}\right| \geqslant \frac{6}{N} \text { for } 1 \leqslant j<k \leqslant q+1
$$

If $\mathbf{m} \in \mathbb{Z}^{q} \backslash\{\mathbf{0}\}, N \geqslant N_{0}(\mathbf{m})$ and $\mathbf{I} \in \mathcal{J}(N)$ we write $\mathcal{E}(\mathbf{m}, N, \mathbf{I})$ for the set of $(E, \mu)$ with the property that if $x_{j} \in E \cap I_{j}$ then

$$
\sum_{j=1}^{q} m_{j} x_{j} \neq 0
$$

By the definition of first category, it suffices to prove the following simpler result.

Lemma 5.4. - If $\mathbf{m} \in \mathbb{Z}^{q+1} \backslash\{\mathbf{0}\}, N \geqslant N_{0}(\mathbf{m})$ and $\mathbf{I} \in \mathcal{J}(N)$, then $\mathcal{E}(\mathbf{m}, N, \mathbf{I})$ is open and dense in $\left(\mathcal{G}, d_{\mathcal{G}}\right)$.

Proof of Theorem 5.3 from Lemma 5.4. - It suffices to show that

$$
\mathcal{E}=\bigcap_{\mathbf{m} \in \mathbb{Z}^{q+1} \backslash\{\mathbf{0}\}} \bigcap_{N \geqslant N_{0}(\mathbf{m})} \bigcap_{\mathbf{I} \in \mathcal{J}(N)} \mathcal{E}(\mathbf{m}, N, \mathbf{I}),
$$

and, since

$$
\mathcal{E} \subseteq \mathcal{E}(\mathbf{m}, N, \mathbf{I})
$$

we need only show that

$$
\mathcal{E} \supseteq \bigcap_{\mathbf{m} \in \mathbb{Z}^{q+1} \backslash\{\mathbf{0}\}} \bigcap_{N \geqslant N_{0}(\mathbf{m})} \bigcap_{\mathbf{I} \in \mathcal{J}(N)} \mathcal{E}(\mathbf{m}, N, \mathbf{I}) .
$$

To this end, let

$$
(E, \mu) \in \bigcap_{\mathbf{m} \in \mathbb{Z}^{q+1} \backslash\{\mathbf{0}\}} \bigcap_{N \geqslant N_{0}(\mathbf{m})} \bigcap_{\mathbf{I} \in \mathcal{J}(N)} \mathcal{E}(\mathbf{m}, N, \mathbf{I})
$$

Suppose $\tilde{\mathbf{m}} \in \mathbb{Z}^{q} \backslash\{\mathbf{0}\}$ and $x_{1}, x_{2}, \ldots x_{q+1}$ are distinct points in $E$. If we choose $\tilde{N} \geqslant N_{0}(\tilde{\mathbf{m}})$ with

$$
\tilde{N} \geqslant 48\left(1+\max _{1 \leqslant i, j \leqslant q+1}\left|x_{i}-x_{j}\right|^{-1}\right)
$$

then we can find $\tilde{\mathbf{I}} \in \mathcal{J}(\tilde{N})$ such that $x_{j} \in \tilde{I}_{j}$. Since

$$
(E, \mu) \in \mathcal{E}(\tilde{\mathbf{m}}, \tilde{N}, \tilde{\mathbf{I}})
$$

we have

$$
\sum_{j=1}^{q+1} m_{j} x_{j} \neq 0
$$

Thus $(E, \mu) \in \mathcal{E}$ and we are done.
We split the proof of Lemma 5.4 into two parts, the easy Lemma 5.5 and the harder Lemma 5.6 (this depends on Lemma 2.5 which we still have to prove).

Lemma 5.5. - If $\mathbf{m} \in \mathbb{Z}^{q+1} \backslash\{\mathbf{0}\}, N \geqslant N_{0}(\mathbf{m})$ and $\mathbf{I} \in \mathcal{J}(N)$, then $\mathcal{E}(\mathbf{m}, N, \mathbf{I})$ is open in $\left(\mathcal{G}, d_{\mathcal{G}}\right)$.

Proof. - Imitate the proof of Lemma 3.6.
Thus the proof of Lemma 5.4 reduces to the proof of the next lemma.
Lemma 5.6. - If $\mathbf{m} \in \mathbb{Z}^{q+1} \backslash\{\mathbf{0}\}, N \geqslant N_{0}(\mathbf{m})$ and $\mathbf{I} \in \mathcal{J}(N)$, then $\mathcal{E}(\mathbf{m}, N, \mathbf{I})$ is dense in $\left(\mathcal{G}, d_{\mathcal{G}}\right)$.

The proof of Lemma 5.6 from Lemma 5.6 will occupy the rest of this section. The next lemma merely serves to establish notation.

Lemma 5.7. - Let $K: \mathbb{R} \rightarrow \mathbb{R}$ be an infinitely differentiable function with the following properties.
(i) $K(x) \geqslant 0$ for all $x \in \mathbb{R}$.
(ii)' $\int_{\mathbb{R}} K(x) d x=1$.
(iii) $K(x)=0$ for $|x| \geqslant 1 / 4$.

If $M$ is a positive integer and we define $K_{M}: \mathbb{T} \rightarrow \mathbb{R}$ by

$$
K_{M}(t)= \begin{cases}M K(M t) & \text { if }|t| \leqslant 1 /(4 M) \\ 0 & \text { otherwise }\end{cases}
$$

then $K_{M}$ is an infinitely differentiable function having the following properties.
(i) $K_{M}(t) \geqslant 0$ for all $t \in \mathbb{T}$.
(ii) $\int_{\mathbb{T}} K_{M}(t) d t=1$.
(iii) $K_{M}(t)=0$ for $|t| \geqslant 1 /(4 M)$.
(iv) $\left|\hat{K}_{M}(r)\right| \leqslant 1$ for all $r$.
(v) There exists a constant $A$, independent of $M$, such that $\left|\hat{K}_{M}(r)\right| \leqslant$ $A(M / r)^{2}$ for all $r \neq 0$.

Proof. - This is entirely straightforward.
Lemma 5.8. - Given $(E, \mu) \in \mathcal{E}$ and $\epsilon>0$, we can find $(F, \tau) \in \mathcal{E}$ such that $d \tau=g d m$ with $g$ infinitely differentiable.

Proof. - Observe that $\left(K_{M} * \mu, E+\left[-M^{-1}, M^{-1}\right]\right) \in \mathcal{E}$ and

$$
d\left(\left(K_{M} * \mu, E+\left[-M^{-1}, M^{-1}\right]\right),(E, \mu)\right) \rightarrow 0
$$

as $M \rightarrow \infty$.
Proof of Lemma 5.6 from Lemma 2.5. - In view of Lemma 5.8, it is sufficient to show that, given $(E, \mu) \in \mathcal{E}$ such that $d \mu=g d m$ with $g$ infinitely differentiable and $\epsilon>0$, we can find $(F, \tau) \in \mathcal{E}(\mathbf{m}, N, \mathbf{I})$ with $d((E, \mu),(F, \tau))<\epsilon$.

Note that, since $g$ is infinitely differentiable there exists a constant $A$ such that

$$
|\hat{g}(r)| \leqslant A|r|^{-3}
$$

for all $r \neq 0$ and, since $(E, \mu) \in \mathcal{E}$, there exists a constant $B$ such that

$$
|\hat{g}(r)| \leqslant B \phi(|r|)
$$

for all $r \neq 0$.
To this end, observe that given $\eta>0$ (to be fixed later), Lemma 2.5 tells us that, since $\eta \psi(|r|) \rightarrow \infty$ as $r \rightarrow \infty$, we can find an infinitely differentiable function $f$ with the following properties.
(i) $f(t) \geqslant 0$ for all $t \in \mathbb{T}$.
(ii) $\hat{f}(0)=1$.
(iii) $|\hat{f}(r)| \leqslant \eta|r|^{-1 /(2 q)}(\log (1+|r|))^{1 / 2} \psi(|r|)=\eta \phi(|r|)$ for all $r \neq 0$.
(iv) If $x_{j} \in \operatorname{supp} f \cap I_{j}$ then

$$
\sum_{j=1}^{q+1} m_{j} x_{j} \neq 0
$$

(v) If $x \in \mathbb{T}$ we can find a $y \in \operatorname{supp} f$ with $|x-y|<\eta$.

Note also that, since $g$ is infinitely differentiable, there exists a constant $A$ such that

$$
|\hat{g}(r)| \leqslant A|r|^{-3}
$$

for all $r \neq 0$. Finally we observe that $r^{-2} \phi(r) \rightarrow 0$ as $r \rightarrow \infty$ so there exists a $C>0$ such that

$$
C \phi(r)>r^{-2}
$$

for all $r \geqslant 1$.
Set $h(t)=g(t) f(t)$ and choose some $1>\delta>0$ (to be fixed later). We seek to estimate $\hat{h}(r)$. If $r \neq 0$

$$
\begin{aligned}
|\hat{h}(r)-\hat{g}(r)| & =\left|\sum_{m \neq 0} \hat{f}(r-m) \hat{g}(m)\right| \leqslant \sum_{m \neq 0}|\hat{f}(r-m)||\hat{g}(m)| \\
& =\sum_{0 \neq|m| \leqslant|r| / 2}|\hat{f}(r-m)||\hat{g}(m)|+\sum_{|m|>|r| / 2}|\hat{f}(r-m)||\hat{g}(m)|
\end{aligned}
$$

Using the remarks about the behaviour of $\phi$ at the beginning of this section, we have

$$
|\hat{f}(r-m)| \leqslant \eta \phi(|r-m|) \leqslant \eta K \phi(r)
$$

whenever $|m| \leqslant|r| / 2$ and

$$
|\hat{f}(r-m)| \leqslant \eta \phi(|r-m|) \leqslant \eta
$$

whenever $|r-m| \neq 0$. Thus

$$
\begin{aligned}
|\hat{h}(r)-\hat{g}(r)| & \leqslant K A \eta \phi(|r|) \sum_{0 \neq|m| \leqslant|r| / 2}|m|^{-3}+A \eta \sum_{|m|>|r| / 2}|m|^{-3} \\
& \leqslant K A \eta \sum_{m \neq 0}|m|^{-3}+A \eta \sum_{|m|>|r| / 2}|m|^{-3} \\
& \leqslant \eta\left(10 K A \phi(|r|)+10 A r^{-2}\right) \leqslant 10 A \eta(K+C) \phi(|r|) \\
& \leqslant \delta \phi(|r|)
\end{aligned}
$$

for all $r \neq 0$ provided only that we choose $\eta$ small enough.

A similar but simpler argument shows that

$$
|\hat{h}(0)-\hat{g}(0)| \leqslant \delta
$$

provided only that we choose $\eta$ small enough. If we now set $H(t)=$ $|\hat{h}(0)|^{-1} h(t)$ then, since

$$
|\hat{H}(r)-\hat{g}(r)| \leqslant|\hat{h}(r)-\hat{g}(r)|+\left|1-\hat{h}(0)^{-1}\right|(|\hat{g}(r)|+|\hat{h}(r)-\hat{g}(r)|)
$$

it follows that, provided we pick $\delta$ small enough,

$$
\sup _{r} \phi(r)^{-1}|\hat{H}(r)-\hat{g}(r)|<\epsilon / 2
$$

Taking $\tau=H d m$ and $F=E \cap \operatorname{supp} f$ we see that $(F, \tau) \in \mathcal{E}$ by construction. Provided $\eta$ is small enough, condition (v) implies that $d_{\mathcal{H}}(E, F)<\epsilon / 2$ and so, using the conclusion of the previous paragraph, $d((E, \mu),(F, \tau))<\epsilon$. Condition (iv) shows that $(F, \tau) \in \mathcal{E}(\mathbf{m}, N, \mathbf{I})$ so we are done.

## 6. Preparations for the main lemma

Before we start the start the proof of Lemma 2.5 in earnest we need to do some cleaning up.

Lemma 6.1. - If $y_{1}, y_{2}, \ldots y_{m}$ are distinct points of $\mathbb{T}, \epsilon>0$ and $\phi$ is as specified at the beginning of section 5, we can find an infinitely differentiable function $f$ with the following properties.
(i) $f(t) \geqslant 0$ for all $t \in \mathbb{T}$.
(ii) $\hat{f}(0)=1$.
(iii) $|\hat{f}(r)| \leqslant \phi(|r|)$ for all $r \neq 0$.
(iv) If $y_{k} \notin \operatorname{supp} f$ for $1 \leqslant k \leqslant m$.
(v) If $x \in \mathbb{T}$ we can find a $y \in \operatorname{supp} f$ with $|x-y|<\epsilon$.

Proof. - (The reader may prefer to supply her own proof.) Choose $K$ in Lemma 5.7 In such a way that

$$
K(x)=\|K\|_{\infty} \text { for }|x| \leqslant 1 / 16
$$

and set

$$
L_{M}(t)=\|K\|_{\infty}^{-1} M^{-1} K_{M}(t)
$$

If we set

$$
g(t)=1-\sum_{j=1}^{m} L_{M}\left(t-y_{k}\right)
$$

then

$$
|\hat{g}(r)| \leqslant \begin{cases}\frac{m\|K\|_{\infty}}{M} & \text { for }|r| \leqslant M \\ \frac{m A\|K\|_{\infty}}{M}\left(\frac{M}{r}\right)^{2} & \text { for }|r| \geqslant M\end{cases}
$$

If $\eta>0$ then, provided only that $M$ is large enough, we have

$$
\frac{m(A+1)\|K\|_{\infty}}{M} \leqslant \eta M^{-3 / 4}
$$

and so

$$
|\hat{g}(r)| \leqslant \eta|r|^{-3 / 4}
$$

for all $r \neq 0$. If we now set $f=\|g\|_{1}^{-1} g$ then, provided that $M$ is large enough, all the conditions of the lemma follow.

Lemma 6.1 gives a proof of Lemma 2.5 in the particular case when all the $m_{j}$ except one are zero. Lemma 2.5 is also trivial in the case when

$$
0 \notin \sum_{j=1}^{q} m_{j} I_{j}
$$

since we can then take $f=1$. Thus we need only prove the following version of Lemma 2.5

Lemma 6.2. - Let $\phi$ be as specified at the beginning of section 5 and let $\epsilon>0$. Suppose $\mathbf{m} \in \mathbb{Z}^{q+1}, m_{1}, m_{2} \neq 1$ and $N$ is a positive integer such that

$$
N \geqslant 12\left(q+1+\sum_{j=1}^{N}\left|m_{j}\right|\right)
$$

Suppose further that we are given $I_{j}=\left[\left(n_{j}-\frac{1}{2}\right) / N,\left(n_{j}+\frac{1}{2}\right) / N\right]$, with $n_{j}$ an integer $[1 \leqslant j \leqslant q+1]$, such that

$$
\left|\frac{n_{j}}{N}-\frac{n_{k}}{N}\right| \geqslant \frac{6}{N} \text { for } 1 \leqslant j<k \leqslant q+1
$$

and

$$
0 \in \sum_{j=1}^{q} m_{j} I_{j}
$$

Then we can find an infinitely differentiable function $f$ with the following properties.
(i) $f(t) \geqslant 0$ for all $t \in \mathbb{T}$.
(ii) $\hat{f}(0)=1$.
(iii) $|\hat{f}(r)| \leqslant \phi(|r|)$ for all $r \neq 0$.
(iv) If $x_{j} \in \operatorname{supp} f \cap I_{j}$ then

$$
\sum_{j=1}^{q+1} m_{j} x_{j} \neq 0
$$

(v) If $x \in \mathbb{T}$ we can find a $y \in \operatorname{supp} f$ with $|x-y|<\epsilon$.

Lemma 6.2 follows in turn from the following result.
Lemma 6.3. - Suppose the hypotheses of Lemma 6.2 hold and in addition we are given infinitely differentiable $g_{j}[1 \leqslant j \leqslant q+1]$ such that
(i) $1 \geqslant g_{j}(t) \geqslant 0$ for all $t \in \mathbb{T}$.
(ii) $g_{j}(t)=0$ if $t \notin\left[\left(n_{j}-\frac{1}{2}-\epsilon\right) / N,\left(n_{j}+\frac{1}{2}+\epsilon\right) / N\right]$.
(iii)' $g_{j}(t)=1$ if $t \in\left[\left(n_{j}-\frac{1}{2}\right) / N,\left(n_{j}+\frac{1}{2}\right) / N\right]$.

Then we can find infinitely differentiable functions $f_{j}$ with the following properties.
(i) ${ }^{\prime \prime} f_{j}(t) \geqslant 0$ for all $t \in \mathbb{T}$.
(ii) ${ }^{\prime \prime} \hat{f}_{j}(0)=\hat{g}_{j}(0)$.
(iii)" $\left|\hat{f}_{j}(r)-\hat{g}_{j}(r)\right| \leqslant(q+1)^{-1} \phi(|r|)$ for all $r \neq 0$.
(iv)" If $x_{j} \in \operatorname{supp} f_{j}$ then

$$
\sum_{j=1}^{q+1} m_{j} x_{j} \neq 0
$$

(v) If $x \in I_{j}$ we can find a $y \in \operatorname{supp} f_{j}$ with $|x-y|<\epsilon$.

Proof of Lemma 6.2 from Lemma 6.3. - Choose $g_{j}$ satisfying conditions (i) $)^{\prime}$, (ii) $)^{\prime}$ and (iii) and set $g=1-\sum_{j=1}^{q+1} g_{j}$. If we choose $f_{j}$ satisfying the conclusions of Lemma 6.3 and set $f=g+\sum_{j=1}^{q+1} f_{j}$, then $f$ satisfies the conclusions of Lemma 6.2.

We can deduce Lemma 6.3 from a result on sums of point masses. Here and elsewhere we write $|E|$ for the number of elements in a finite set $E$.

Lemma 6.4. - Suppose the hypotheses of Lemma 6.3 hold. Then we can find $N_{0}, N_{1}, B \geqslant 1$ and $\gamma>0$ with the following properties. If $n \geqslant N_{1}$ we can find finite sets of points $E_{j}[1 \leqslant j \leqslant q+1]$ such that writing

$$
\mu_{j}=\left|E_{j}\right|^{-1}\left\|g_{j}\right\|_{1} \sum_{x \in E_{j}} \delta_{x}
$$

the following conditions hold.
(1) $\left|\hat{\mu}_{j}(r)-\hat{g}_{j}(r)\right| \leqslant 2^{-1}(q+1)^{-1} \phi(|r|)$ for all $|r| \leqslant N_{0}$.
(2) $\left|\hat{\mu}_{j}(r)\right|+\left|\hat{g}_{j}(r)\right| \leqslant(q+1)^{-1} \phi(|r|)$ for all $N_{0} \leqslant|r| \leqslant n^{2(q+1)}$.
(3) If $x_{j} \in E_{j}+\left[-4 \gamma n^{-q}, 4 \gamma n^{-q}\right]$ then

$$
\sum_{j=1}^{q+1} m_{j} x_{j} \neq 0
$$

(4) If $x \in I_{j}$ we can find a $y \in E_{j}$ with $|x-y|<\epsilon$.

Proof of Lemma 6.3 from Lemma 6.4. - let $M=M(n)$ be the integer satisfying

$$
\left(\gamma n^{-q}\right)^{-1}+1 \geqslant M>\left(\gamma n^{-q}\right)^{-1}
$$

and set $f_{j}=\mu_{j} * K_{M}$ where $K_{M}$ is defined as in Lemma 5.7. Automatically the $f_{j}$ satisfy the conclusions of Lemma 6.3. with the possible exception of (iii)". We now show that (iii)" holds, provided only that $n$ is large enough.

First observe that, provided only that $n$ is large enough, standard results on approximate identities tell us that

$$
\left|\hat{f}_{j}(r)-\hat{\mu}_{j}(r)\right| \leqslant 2^{-1}(q+1)^{-1} \phi(|r|)
$$

and so, using (1),

$$
\left|\hat{g}_{j}(r)-\hat{f}_{j}(r)\right| \leqslant(q+1)^{-1} \phi(|r|)
$$

for all $|r| \leqslant N_{0}$ provided only that $n$ is large enough. Next we note that

$$
\left|\hat{f}_{j}(r)\right|=\left|\hat{\mu}_{j}(r)\right|\left|\hat{K}_{M}(r)\right| \leqslant\left|\hat{\mu}_{j}(r)\right|
$$

so, using (2),

$$
\left|\hat{f}_{j}(r)\right|+\left|\hat{g}_{j}(r)\right| \leqslant(q+1)^{-1} \phi(|r|)
$$

whence

$$
\left|\hat{\mu}_{j}(r)-\hat{f}_{j}(r)\right| \leqslant(q+1)^{-1} \phi(|r|)
$$

for all $N_{0} \leqslant|r| \leqslant n^{2(q+1)}$.
Note that, since $g_{j}$ is infinitely differentiable, there exists a constant $C$ such that

$$
\left|\hat{g}_{j}(r)\right| \leqslant C|r|^{-1}
$$

for all $r \neq 0$. Thus, provided only that $n$ is large enough,

$$
\left|\hat{g}_{j}(r)\right| \leqslant(q+1)^{-1} \phi(|r|) / 2
$$

for all $n^{2 q}+1 \leqslant|r|$. Using the equality $\left|\hat{f}_{j}(r)\right|=\left|\hat{\mu}_{j}(r)\right|\left|\hat{K}_{M}(r)\right|$, we observe that

$$
\left|\hat{f}_{j}(r)\right| \leqslant A(M / r)^{2} \leqslant 2 A \gamma^{-2}\left(n^{q} / r\right)^{2} \leqslant 2 A \gamma^{-2}|r|^{-1}
$$

for $n^{2 q}+1 \leqslant|r|$. Thus, provided only that $n$ is large enough,

$$
\left|\hat{f}_{j}(r)\right| \leqslant(q+1)^{-1} \phi(|r|) / 2
$$

and so

$$
\left|\hat{f}_{j}(r)-\hat{g}_{j}(r)\right| \leqslant(q+1)^{-1} \phi(|r|)
$$

for all $n^{2 q}+1 \leqslant|r|$ and so we are done.

## 7. Completion of the proof of the main lemma

In this final section we obtain Lemma 6.4 by means of a probabilistic construction. All parts of following theorem are well known (see for example [1]) but it may be helpful to recall the proofs.

Theorem 7.1. - (i) If $Y$ is a real valued random variable with $|Y| \leqslant$
1 and $\mathbb{E} Y=0$ then

$$
\mathbb{E} e^{\lambda Y} \leqslant e^{\lambda^{2}}
$$

(ii) If $Y_{1}, Y_{2} \ldots$ are independent real valued random variable with $\left|Y_{k}\right| \leqslant 1$ and $\mathbb{E} Y_{k}=0$ then

$$
\operatorname{Pr}\left(\sum_{k=1}^{n} Y_{k} \geqslant y\right) \leqslant e^{-y^{2} / 4 n}
$$

(iii) If $Z_{1}, Z_{2} \ldots$ are independent complex valued random variable with $\left|Z_{k}\right| \leqslant 1$ and $\mathbb{E} Z_{k}=0$ then

$$
\operatorname{Pr}\left(\left|\sum_{k=1}^{n} Z_{k}\right| \geqslant y\right) \leqslant 4 e^{-y^{2} / 4 n}
$$

(iv) Suppose $U_{1}, U_{2} \ldots$ are independent identically distributed random variables taking values on $\mathbb{T}$. If

$$
\operatorname{Pr}\left(U_{1} \in[a, b)\right)=\mu([a, b))
$$

for some probability measure $\mu$ then

$$
\operatorname{Pr}\left(\left|n^{-1} \sum_{j=1}^{n} e^{i r U_{k}}-\hat{\mu}(r)\right| \geqslant y\right) \leqslant 4 e^{-y^{2} /(16 n)}
$$

(v) Suppose $0 \leqslant \alpha \leqslant 1$ and $W_{1}, W_{2}, \ldots$ are independent complex valued random variables with $\left|W_{k}\right| \leqslant 1$ and $\left|\mathbb{E} W_{k}\right| \leqslant \alpha$ is as in (iv). Then

$$
\operatorname{Pr}\left(\left|n^{-1} \sum_{j=1}^{n} W_{j}\right| \geqslant \alpha+y\right) \leqslant 4 e^{-y^{2} /(16 n)}
$$

Proof. -
(i) The result is immediate if $|\lambda| \geqslant 1$. If $|\lambda| \leqslant 1$,
$\mathbb{E} e^{\lambda Y}=\sum_{r=0}^{\infty} \mathbb{E} Y^{r} \frac{\lambda^{r}}{r!}=1+\sum_{r=2}^{\infty} \mathbb{E} Y^{r} \frac{\lambda^{r}}{r!} \leqslant=1+\sum_{r=2}^{\infty} \frac{|\lambda|^{r}}{r!} \leqslant \sum_{r=0}^{\infty} \frac{|\lambda|^{2 r}}{r!}=e^{\lambda^{2}}$.
(ii) Observe that the random variables $e^{\lambda Y_{k}}$ are independent so

$$
\mathbb{E} e^{\lambda \sum_{j=1}^{n} Y_{k}}=\mathbb{E} \prod_{j=1}^{n} e^{\lambda Y_{k}}=\prod_{j=1}^{n} \mathbb{E} e^{\lambda Y_{k}} \leqslant e^{n \lambda^{2}}
$$

Thus by a Tchebychev estimate

$$
\operatorname{Pr}\left(\sum_{k=1}^{n} Y_{k} \geqslant y\right) \leqslant e^{-\lambda y} \mathbb{E} e^{\lambda \sum_{j=1}^{n} Y_{k}}=e^{n \lambda^{2}-\lambda y}
$$

and setting $\lambda=y / 2 n$ we have the desired result.
(iii) Apply part (ii) to $\Re Z_{k},-\Re Z_{k}, \Im Z_{k}$ and $-\Im Z_{k}$.
(iv) Observe that

$$
\mathbb{E} e^{i r U_{k}}=\hat{\mu}(r),
$$

so applying part (iii) with $Z_{k}=\left(e^{i r U_{k}}-\hat{\mu}(r)\right) / 2$ gives the required result.
(v) Observe that

$$
\begin{aligned}
\operatorname{Pr}\left(\left|n^{-1} \sum_{k=1}^{n} W_{k}\right| \geqslant|\hat{\mu}(r)|+y\right) & \leqslant \operatorname{Pr}\left(\left|n^{-1} \sum_{k=1}^{n} W_{k}\right| \geqslant\left|\mathbb{E} W_{1}\right|+y\right) \\
& \leqslant \operatorname{Pr}\left(\left|n^{-1} \sum_{k=1}^{n}\left(W_{k}-\mathbb{E} W_{k}\right)\right| \geqslant y\right) \\
& \leqslant 4 e^{-y^{2} /(16 n)}
\end{aligned}
$$

as in part (iv).

We now state our probabilistic version of Lemma 6.4.
Lemma 7.2. - Suppose the hypotheses of Lemma 6.3 hold. Set $M=$ $\sum_{j=1}^{q+1}\left|m_{j}\right|$. Then we can find $N_{0}, N_{1}, B \geqslant 1$ and $\gamma>0$ with the following properties such that whenever $n \geqslant N_{1}$ the following is true.

Suppose $X_{j k}$ are independent random variables taking values on $\mathbb{T}[1 \leqslant$ $j \leqslant q+1,1 \leqslant k \leqslant n]$ such that $X_{j k}$ has probability density $\left\|g_{j}\right\|_{1}^{-1} g_{j}$. If $2 \leqslant j \leqslant q+1$ take

$$
E_{j}=\left\{X_{j k}: 1 \leqslant k \leqslant n\right\}
$$

set

$$
\tilde{E}_{1}=\left\{X_{j k}: 1 \leqslant k \leqslant n\right\}
$$

and

$$
E_{1}=\left\{x \in E_{1}: 0 \notin\left[-4 M \gamma n^{-q}, 4 M \gamma n^{-q}\right]+m_{1} x+\sum_{j=2}^{q+1} m_{j} E_{j}\right\}
$$

If we take

$$
\mu_{j}=\left|E_{j}\right|^{-1}\left\|g_{j}\right\|_{1} \sum_{x \in E_{j}} \delta_{x}
$$

then, with probability at least $1 / 2$, following conditions hold.
(1) $\left|\hat{\mu}_{j}(r)-\hat{g}_{j}(r)\right| \leqslant 2^{-1}(q+1)^{-1} \phi(|r|)$ for all $|r| \leqslant N_{0}$.
(2) $\left|\hat{\mu}_{j}(r)\right|+\left|\hat{g}_{j}(r)\right| \leqslant(q+1)^{-1} \phi(|r|)$ for all $N_{0} \leqslant|r| \leqslant n^{2(q+1)}$.
(3) If $x_{j} \in E_{j}+\left[-4 \gamma n^{-q}, 4 \gamma n^{-q}\right]$, then

$$
\sum_{j=1}^{q+1} m_{j} x_{j} \neq 0
$$

(4) If $x \in I_{j}$ we can find a $y \in E_{j}$ with $|x-y|<\epsilon$.

Since any event which has strictly positive probability must have an instance Lemma 7.2 follows from Lemma 6.4.

Most of Lemma 7.2 is easy to prove.
Lemma 7.3. - Suppose the hypotheses of Lemma 6.3 hold. Set $M=$ $\sum_{j=1}^{q+1}\left|m_{j}\right|$. Then we can find $N_{0}^{\prime}, N_{1}^{\prime}$ and $B^{\prime} \geqslant 1$ such that whenever $n \geqslant N_{1}^{\prime}$ and $B \geqslant B^{\prime}$ the following is true.

Suppose $X_{j k}$ are independent random variables taking values on $\mathbb{T}[1 \leqslant$ $j \leqslant q+1,1 \leqslant k \leqslant n]$ such that $X_{j k}$ has probability density $\left\|g_{j}\right\|_{1}^{-1} g_{j}$ and suppose $\gamma>0$. If $2 \leqslant j \leqslant q+1$ take

$$
E_{j}=\left\{X_{j k}: 1 \leqslant k \leqslant n\right\}
$$

set

$$
\tilde{E}_{1}=\left\{X_{j k}: 1 \leqslant k \leqslant n\right\}
$$

and

$$
E_{1}=\left\{x \in E_{1}: 0 \notin\left[-4 M \gamma n^{-q}, 4 M \gamma n^{-q}\right]+m_{1} x+\sum_{j=2}^{q+1} m_{j} E_{j}\right\}
$$

If we take

$$
\mu_{j}=\left|E_{j}\right|^{-1}\left\|g_{j}\right\|_{1} \sum_{x \in E_{j}} \delta_{x}
$$

then, with probability at least $3 / 4$, the following conditions hold.
(1) $)^{\prime}\left|\hat{\mu}_{j}(r)-\hat{g}_{j}(r)\right| \leqslant 2^{-1}(q+1)^{-1} \phi(|r|)$ for all $|r| \leqslant N_{0}^{\prime}$ and $2 \leqslant j \leqslant q+1$.
(2)' $\left|\hat{\mu}_{j}(r)\right|+\left|\hat{g}_{j}(r)\right| \leqslant(q+1)^{-1} \phi(|r|)$ for all $N_{0} \leqslant|r| \leqslant n^{2(q+1)}$ and $2 \leqslant j \leqslant q+1$.
(3) If $x_{j} \in E_{j}+\left[-4 \gamma n^{-q}, 4 \gamma n^{-q}\right]$, then

$$
\sum_{j=1}^{q+1} m_{j} x_{j} \neq 0
$$

(4) If $2 \leqslant j \leqslant q+1$ and $x \in I_{j}$ we can find a $y \in E_{j}$ with $|x-y|<\epsilon$.

Proof. - Observe that (3) is always true by virtue of the definition of $E_{1}$. The weak law of large numbers tells us that, provided only that $n$ is large enough, condition $(4)^{\prime}$ will hold with probability at least $7 / 8$.

Since $g_{j}$ is once continuously differentiable we can find a $C_{j}$ such that

$$
\left|\hat{g}_{j}(r)\right| \leqslant C_{j}|r|^{-1}
$$

for all $|r|>0$ and so we can find an $N_{0}^{\prime}$ such that

$$
\left|\hat{g}_{j}(r)\right| \leqslant 4^{-1}(q+1)^{-1} \phi(|r|)
$$

for all $|r| \leqslant N_{0}^{\prime}$ and $2 \leqslant j \leqslant q+1$.
By Theorem 7.1

$$
\begin{aligned}
\operatorname{Pr}\left(\left|\hat{\mu}_{j}(r)-\hat{g}_{j}(r)\right|\right. & \left.\geqslant B n^{-1 / 2}(\log n)^{1 / 2}\right) \\
& =\operatorname{Pr}\left(\left|n^{-1} \sum_{k=1}^{n} e^{i r X_{k}}-\hat{\mu}(r)\right| \geqslant B n^{-1 / 2}(\log n)^{1 / 2}\right) \\
& \leqslant 4 e^{-B \log n /(16 n)} .
\end{aligned}
$$

Thus, if we choose $B \geqslant 64(q+1)$, we have

$$
\operatorname{Pr}\left(\left|\hat{\mu}_{j}(r)-\hat{g}_{j}(r)\right| \geqslant B n^{-1 / 2}(\log n)^{1 / 2}\right) \leqslant 4 n^{-4(q+1)}
$$

for all $r$ and all $2 \leqslant j \leqslant q+1$. Thus provided only that $n$ is large enough,

$$
\operatorname{Pr}\left(\left|\hat{\mu}_{j}(r)-\hat{g}_{j}(r)\right| \geqslant B n^{-1 / 2}(\log n)^{1 / 2}\right) \leqslant 4 n^{-4(q+1)}
$$

will hold with probability at least $7 / 8$.
Using the results of the two previous paragraphs we see that conditions $(1)^{\prime}$ and $(2)^{\prime}$ will both hold (with probability at least $7 / 8$ ) provided only that $n$ is large enough. The result follows.

We now prove the harder part of Lemma 7.2.
Lemma 7.4. - Suppose the hypotheses of Lemma 6.3 hold. Set $M=$ $\sum_{j=1}^{q+1}\left|m_{j}\right|$. Then we can find $N_{0}^{\prime \prime}, N_{1}^{\prime \prime}, B^{\prime \prime} \geqslant 1$ and $\gamma>0$ such that whenever $n \geqslant N_{1}^{\prime \prime}$ and $B \geqslant B_{1}$ the following is true.

Suppose $X_{j k}$ are independent random variables taking values on $\mathbb{T}[1 \leqslant$ $j \leqslant q+1,1 \leqslant k \leqslant n]$ such that $X_{j k}$ has probability density $\left\|g_{j}\right\|_{1}^{-1} g_{j}$. If $2 \leqslant j \leqslant q+1$ take

$$
E_{j}=\left\{X_{j k}: 1 \leqslant k \leqslant n\right\}
$$

set

$$
\tilde{E}_{1}=\left\{X_{j k}: 1 \leqslant k \leqslant n\right\}
$$

and

$$
E_{1}=\left\{x \in \tilde{E}_{1}: 0 \notin\left[-4 M \gamma n^{-q}, 4 M \gamma n^{-q}\right]+m_{1} x+\sum_{j=2}^{q+1} m_{j} E_{j}\right\}
$$

If we take

$$
\mu_{1}=\left|E_{1}\right|^{-1}\left\|g_{j}\right\|_{1} \sum_{x \in E_{1}} \delta_{x}
$$

then, with probability at least $3 / 4$, following conditions hold.
$(1)^{\prime \prime}\left|\hat{\mu}_{1}(r)-\hat{g}_{1}(r)\right| \leqslant 2^{-1}(q+1)^{-1} \phi(|r|)$ for all $|r| \leqslant N_{0}^{\prime \prime}$.
(2) $)^{\prime \prime}\left|\hat{\mu}_{1}(r)\right|+\left|\hat{g}_{1}(r)\right| \leqslant(q+1)^{-1} \phi(|r|)$ for all $N_{0}^{\prime \prime} \leqslant|r| \leqslant n^{2(q+1)}$.
(4)" If $x \in I_{1}$ we can find a $y \in E_{1}$ with $|x-y|<\epsilon$.

Proof. - Let

$$
\begin{aligned}
E_{*} & =E_{1} \backslash \tilde{E}_{1} \\
& =\left\{x \in \tilde{E}_{1}: 0 \in\left[-4 M \gamma n^{-q}, 4 M \gamma n^{-q}\right]+m_{1} x+\sum_{j=2}^{q+1} m_{j} E_{j}\right\}
\end{aligned}
$$

and

$$
\tau=n^{-1}\|g\|_{1} \sum_{x \in E_{*}} \delta_{x}
$$

The main step of the proof involves finding an upper bound for $\hat{\tau}(r)$ which holds with high probability independent of the choice of $\gamma$.

First observe that, if we set $W_{k}=e^{i r X_{1 k}}$ when $X_{1 k} \in E_{*}$ and $W_{k}=0$ otherwise, then

$$
\hat{\tau}(r)=\|g\|_{1} \sum_{k=1}^{n} Z_{k}
$$

the $W_{k}$ satisfy the conditions of Theorem 7.1 (v).
Since $X_{2 k}$ has density function $g_{2} /\left\|g_{2}\right\|_{1}$ it follows that, to first order in $\delta t$

$$
\operatorname{Pr}\left(m_{2} X_{2 k} \in[t, t+\delta t]\right)=G(t) \delta t
$$

where $G$ is differentiable density function with first and second derivatives bounded by some $K_{1}$ depending only on $m_{2}$ and $g_{2}$. Thus

$$
\left.\operatorname{Pr}\left(m_{2} X_{2 k} \in[t, t+\delta t]: \text { for some } 1 \leqslant k \leqslant n\right)\right\}=n G(t) \delta t
$$

to first order in $\delta t$ and

$$
\operatorname{Pr}\left(\left(m_{2} E_{2}+m_{3} E_{3}+\cdots+m_{q+1} E_{q+1}\right) \cap[t, t+\delta t] \neq \emptyset\right)=n G * H(t) \delta t
$$

for some $H$. We observe that $G * H$ is differentiable density function with first and second derivatives bounded by $K_{1}$. It follows that, if $t$ is fixed
$\operatorname{Pr}\left(\left(t+m_{2} E_{2}+m_{3} E_{3}+\cdots+m_{q+1} E_{q+1}\right) \cap\left[-4 M \gamma n^{-q}, 4 M \gamma n^{-q}\right] \neq \emptyset\right)=F_{\gamma}(t)$
where $F_{\gamma}$ has continuous first and second derivatives bounded by $\|F\|_{\gamma}^{-1} K_{1}$. Thus the density function $G_{\gamma}$ of $X_{11}$ given that

$$
\left(m_{1} X_{11}+m_{2} E_{2}+m_{3} E_{3}+\cdots+m_{q+1} E_{q+1}\right) \cap\left[-4 M \gamma n^{-q}, 4 M \gamma n^{-q}\right] \neq \emptyset
$$

has a continuous derivative bounded by $K_{2}$, where $K_{2}$ is independent both of $\gamma$ and $n$.

We now have

$$
\left|\mathbb{E} W_{k}\right|=\operatorname{Pr}\left(W_{k} \neq 0\right)\left|\mathbb{E}\left(W_{k} \mid W_{k} \neq 0\right)\right|=\left|\hat{G}_{\gamma}(r)\right| \leqslant \frac{K_{2}}{|r|}
$$

for $r \neq 0$. Using Theorem 7.1 (v), we see that that, if take $B \geqslant 64(q+1)$ then provided we take $n$ large enough, there is a probability at least $31 / 32$ that

$$
\begin{equation*}
|\hat{\tau}(r)| \leqslant K_{2}|r|^{-1}+B n^{-1 / 2}(\log n)^{1 / 2} \tag{1}
\end{equation*}
$$

for all $1 \leqslant|r| \leqslant n^{2(q+1)}$.
Next we observe that

$$
\operatorname{Pr}\left(X_{1 k} \in E_{*}\right) \leqslant n^{q} \times\left(8 M \gamma n^{-q}\right)=8 M \gamma
$$

so the expected number of points in $E^{*}$ is no greater than $8 M \gamma$. Since

$$
y \operatorname{Pr}(Y \geqslant y) \leqslant \mathbb{E} Y
$$

it follows that given $\eta>0$ (to be fixed later) we can choose $\gamma$ so small that with probability at least $31 / 32 E^{*}$ contains at most $\eta n$ points and so

$$
\begin{equation*}
\|\tau\| \leqslant \eta\|g\|_{1} \tag{2}
\end{equation*}
$$

If we set

$$
\mu=n^{-1}\|g\|_{1} \sum_{x \in E_{1}} \delta_{x}
$$

the argument of Lemma 7.3 shows that, provided that $n$ is large enough, then with probability at least $31 / 32$,

$$
\begin{equation*}
\left|\hat{\mu}(r)-\hat{g}_{1}(r)\right| \leqslant B n^{-1 / 2}(\log n)^{1 / 2} \tag{3}
\end{equation*}
$$

for all $|r| \leqslant n^{2(q+1)}$. Since $g_{1}$ is continuously differentiable there exists a $C$ such that $\left|\hat{g}_{1}(r)\right| \leqslant C|r|^{-1}$ for $r \neq 0$.

For the moment we suppose simply that $\eta \leqslant 1 / 2$. Since

$$
\mu_{1}=\left(\|g\|_{1}-\|\tau\|\right)^{-1}(\mu-\tau)
$$

it follows that, if $(\mathbf{1}),(\mathbf{2})$ and (3) hold

$$
\begin{aligned}
\left|\hat{\mu}_{1}(r)\right|+\left|\hat{g}_{1}(r)\right| & \leqslant 2(|\hat{\mu}(r)|+\mid \hat{\tau}(r))+\left|\hat{g}_{1}(r)\right| \\
& \leqslant 2\left(\left|\hat{\mu}(r)-\hat{g}_{1}(r)\right|+\mid \hat{\tau}(r)\right)+3\left|\hat{g}_{1}(r)\right| \\
& \leqslant 4 B n^{-1 / 2}(\log n)^{1 / 2}+\frac{2 K_{2}+C}{|r|}
\end{aligned}
$$

for all $|r| \leqslant n^{2(q+1)}$. Thus we can find $N_{0}^{\prime \prime}$ independent of $\eta$ (provided $\eta<1 / 2)$ such that, if (1), (2) and (3), hold

$$
\left|\hat{\mu}_{1}(r)\right|+\left|\hat{g}_{1}(r)\right| \leqslant 4^{-1}(q+1)^{-1} \phi(|r|)
$$

for all $|r| \geqslant N_{0}^{\prime \prime}$.
Once $N_{0}^{\prime \prime}$ is fixed, we see that, provided only that $\eta$ (and so $\gamma$ ) is taken sufficiently small, we will have

$$
\left|\hat{\mu}_{1}(r)-\hat{\mu}(r)\right| \leqslant B n^{-1 / 2}(\log n)^{1 / 2}
$$

for all $|r| \leqslant N_{0}^{\prime}$ and so

$$
\left|\hat{\mu}_{1}(r)-\hat{g}_{1}(r)\right| \leqslant 2^{-1}(q+1)^{-1} \phi(|r|)
$$

for all $|r| \leqslant N_{0}^{\prime \prime}$ whenever (2) (3) hold and $n$ is sufficiently large.
Once $\gamma$ is fixed, the weak law of large numbers tells us that, provided only that $n$ is large enough, condition (4) " will hold with probability at least $31 / 32$. Thus, provided only that $n$ is large enough (1), (2), (3) and (4) ${ }^{\prime \prime}$ will hold simultaneously with probability at least $7 / 8$ and imply the conclusions of the lemma.

## BIBLIOGRAPHY

[1] J. P. Kahane, Some random series of functions, Second edition. Cambridge Studies in Advanced Mathematics, no. 5, Cambridge University Press, Cambridge, 1985.
[2] R. Kaufman, "A functional method for linear sets", Israel J. Math. 5 (1967), p. 185187.
[3] , "Small subsets of finite Abelian groups", Annales de l'Institut Fourier 18 (1968), p. 99-102.
[4] T. W. K. Körner, "Measures on independent sets, a quantitative version of a theorem of Rudin", Proc. Amer. Math. Soc. 135 (2007), no. 12, p. 3823-3832.
[5] K. Kuratowski, Topology, vol. I, Academic Press, New York-London, Państwowe Wydawnictwo Naukowe, Warsaw, 1966 (Translated from the French by J. Jaworowski).

Manuscrit reçu le 10 janvier 2008, accepté le 21 mars 2008.

Thomas W. KÖRNER
DPMMS
Centre for Mathematical Sciences
Clarkson Road
Cambridge (England)
twk@dpmms.cam.ac.uk

