# Some results on convolution idempotents 

P Charantej Reddy<br>IIT Hyderabad<br>Email: ee18resch01010@iith.ac.in

Aditya Siripuram<br>IIT Hyderabad<br>Email: staditya@iith.ac.in

Brad Osgood<br>Stanford University<br>Email: osgood@stanford.edu


#### Abstract

We consider the problem of recovering $N$ length vectors $h$ that vanish on a given set of indices and satisfy $h * h=h$. We give some results on the structure of such $h$ when $N$ is a product of two primes, and investigate some bounds and their connections to certain graphs defined on $\mathbb{Z}_{N}$.


Index Terms-Discrete Fourier transform, sampling, Ramanujan's sums, Fuglede's conjecture, convolution, idempotents

## I. Introduction

We say that an $N$ length vector $h$ is a convolution idempotent if $h * h=h$, where $*$ denotes discrete circular convolution. We consider the problem of recovering $h$ when some of its entries are known to be zero. This problem has applications to sampling analog signals, and is possibly a useful connection to make progress on the Fuglede conjecture [1]. In this paper, we briefly review the motivations for considering such a problem. We then present some results when $N$ is a product of two primes. This is a followup to our work in [2] where we characterized all possible idempotents with given zero sets, in the case when $N$ is a prime power.

Some preliminary notations and observations before we proceed. We denote by $\mathbb{Z}_{N}$ the set of integers modulo $N$, and by $\mathcal{F}_{N}$ the $N \times N$ Discrete Fourier transform (DFT) matrix (we omit $N$ if it is apparent from the context) [3], [4], so that for any $x \in \mathbb{C}^{N}$, we have

$$
\mathcal{F} x(n)=\sum_{j \in \mathbb{Z}_{N}} x(j) e^{-2 \pi i n j / N}, \text { for } n \in \mathbb{Z}_{N}
$$

The discrete (circular) convolution of two signals $x, y \in \mathbb{C}^{N}$ is

$$
(x * y)(n)=\sum_{j \in \mathbb{Z}_{N}} x(j) y(n-j), \text { for } n \in \mathbb{Z}_{N}
$$

Recall that $\mathcal{F}(x * y)=(\mathcal{F} x)(\mathcal{F} y)$, so that for a convolution idempotent $h$ we have $(\mathcal{F} h)^{2}=\mathcal{F} h$. Thus the entries of $\mathcal{F} h$ are either 0 or 1 , and an idempotent $h$ is equivalently defined by the support of $\mathcal{F} h$. That is, for a (support) set $\mathcal{J} \subseteq \mathbb{Z}_{N}$, with $1_{\mathcal{J}} \in \mathbb{C}^{N}$ denoting the indicator of $\mathcal{J}$ the idempotent corresponding to $\mathcal{J}$ is $h_{\mathcal{J}}=\mathcal{F}^{-1} 1_{\mathcal{J}}$.

Just as an idempotent is characterized by its support $\mathcal{J} \subseteq$ $\mathbb{Z}_{N}$, a foundational result is that the zero set of an idempotent is characterized by a subset of divisors of $N$ (stated as Lemma 1, below). For this, we introduce the following notation: For a vector $h \in \mathbb{C}^{N}$, let $\mathcal{Z}(h) \subseteq \mathbb{Z}_{N}$ denote the indices where $h$ vanishes. Let $\mathcal{D}_{N}$ be the set of all divisors of $N$ in $\mathbb{Z}_{N}$ (so 978-1-7281-6432-8/20/\$31.00 ©2020 IEEE
omitting $N$ ), let $(i, N)$ denote the greatest common divisor $(\operatorname{gcd})$ of $i$ and $N$, and let

$$
\begin{equation*}
\mathscr{A}_{N}(k)=\left\{i \in \mathbb{Z}_{N}:(i, N)=k\right\} . \tag{1}
\end{equation*}
$$

Lemma 1: The zero-set $\mathcal{Z}(h)$ is the disjoint union

$$
\mathcal{Z}(h)=\left\{i \in \mathbb{Z}_{N}:(i, N) \in \mathcal{D}(h)\right\}=\bigcup_{k \in \mathcal{D}(h)} \mathscr{A}_{N}(k)
$$

for some set of divisors $\mathcal{D}(h) \subseteq \mathcal{D}_{N}$.
We call $\mathcal{D}(h)$ the zero-set divisors of $h$. The lemma appears in many different forms and contexts, see [5], [6], [7] or [8, Theorem 2.1] for example: the key ingredient in the proof is the structure of cyclotomic polynomials.

For example if $N=8, \mathcal{J}=\{0,1,4,5\}$ then $\mathcal{Z}\left(h_{\mathcal{J}}\right)=$ $\{1,3,4,5,7\}=\{1,3,5,7\} \cup\{4\}=\mathscr{A}_{8}(1) \cup \mathscr{A}_{8}(4)$; here $\mathcal{D}\left(h_{\mathcal{J}}\right)=\{1,4\}$.

We can ask if a converse of Lemma 1 holds; i.e. given a set $\mathcal{Z}$ of the form in Lemma 1 (i.e. $\mathcal{Z}$ is a union of some $\mathscr{A}_{N}(k)$ ), is there an idempotent $h$ whose zero set is $Z$ ? We formulate this as:

Problem $\mathfrak{i}_{N}(\mathcal{D})$ : Given a positive integer $N$ and a set of divisors $\mathcal{D} \subseteq \mathcal{D}_{N}$ let

$$
\begin{equation*}
\mathcal{Z}=\left\{i \in \mathbb{Z}_{N}:(i, N) \in \mathcal{D}\right\} \tag{2}
\end{equation*}
$$

Find all index sets $\mathcal{J}$ such that the idempotent $h_{\mathcal{J}}=$ $\mathcal{F}^{-1} 1_{\mathcal{J}}$ vanishes on $\mathcal{Z}$.
Note that this is slightly different from a converse of Lemma 1: For $\mathcal{J}$ to be a solution to $\mathfrak{i}(\mathcal{D})$ we only ask if $h_{\mathcal{J}}$ vanishes on $\mathcal{Z}$. In particular, $h$ may vanish outside $\mathcal{Z}$ (i.e. have a bigger zero set) as well.

In our previous work [6], we motivated the zero set problem $\mathfrak{i}(\mathcal{D})$ in the context of sampling. We revisit and summarize this in Section II. In [6], we also gave a complete characterization of all solutions to $\mathfrak{i}(\mathcal{D})$ when $N=p^{M}$ is a prime power, using base $-p$ expansions of elements of $\mathcal{J}$. In this work, we build a case for solving $\mathfrak{i}(\mathcal{D})$ when $N$ has more than one prime factor, and generalize some of the results of [2]. The main contribution here is to characterize all solutions to $\mathfrak{i}(\mathcal{D})$ when $N=p q$ and $\mathcal{D}=\{1\},\{p, q\},\{1, p\}$ or $\{1, q\}$ (Theorem 1 , Corollary 1, 2). Theorem 1 seems connected to results on vanishing sums of roots of unity [9], though the techniques we employ are different. The proof of Theorem 1 uses properties of Ramanujan sums [10], which were recently used for signal processing [11], [12]. In section IV we investigate some bounds on the smallest solution to $\mathfrak{i}(\mathcal{D})$, and relate it to the minrank [13] of certain graphs introduced in [6].

## II. Motivation

## A. Sampling

We describe a problem in the traditional multicoset sampling setting [14], [15] in which the zero set problem naturally arises. See [2] for details: here we will give a concise overview of the key ideas.

We are interested in sampling signals that have a fragmented spectrum: i.e. for any signal $f$ in the space, the Fourier transform $\mathcal{F} f(s)$ is non zero only when the frequency $s$ is in $\cup_{n \in \mathfrak{F}}[n, n+1]$. See Figure 1 for an example of a signal in such a space, with $\mathfrak{F}=\{0,2\}$.


Fig. 1: Example signal with two fragments, for $\mathfrak{F}=\{0,2\}$.

If we sample the signal from Fig 1 at the Nyquist rate, we need to take at least 3 samples per second. Instead, consider sampling the signal with the sampling pattern shown in Fig 2:


Fig. 2: Example sampling pattern in the case $\mathfrak{F}=\{0,2\}$. Samples are taken at every second, and at a 0.25 second offset

Note that with the sampling pattern of Figure 2, we take on average 2 samples per second. With elementary Fourier analysis, one can show that the sampled signal has a spectrum shown in Fig 3 (see [2] for details, the main idea is from [14], [15]).


Fig. 3: Spectrum of signal in Fig 1 sampled with the spectrum from Fig 2.

The crucial observation from Fig 3 is that the original signal fragments are intact: aliasing primary occurs in the islands
where the signal was non-existent anyway. This idea can be generalized to signal spaces with arbitrary locations of the fragments $\mathfrak{F}$, by picking a sampling pattern of the form

$$
\begin{equation*}
p_{\mathcal{J}}(t)=\sum_{m \in \mathcal{J}} \delta(t-m / N) \tag{3}
\end{equation*}
$$

where $\mathcal{J}$ and $N$ are parameters that need to be picked. For example, for the sampling pattern in Fig 1, we have $N=$ $4, \mathcal{J}=\{0,1\}$. Now we let $\mathcal{Z}$ be the difference set of $\mathfrak{F}$, and $\mathcal{D}$ the corresponding gcds, i.e.,

$$
\begin{gathered}
\mathcal{Z}=\left\{k_{1}-k_{2}: k_{1}, k_{2} \in \mathfrak{F}\right\} \\
\mathcal{D}=\{d:(k, N)=d \text { for some } \mathrm{k} \in \mathfrak{F} .\}
\end{gathered}
$$

We can show ( [2, Proposition 2]) that any $\mathcal{J}$ providing a solution to the zero set problem $\mathfrak{i}(\mathcal{D})$ will ensure that the signal can be recovered from the samples taken according to (3). The average sampling rate would then be $|\mathcal{J}|$, which could potentially be much smaller than the Nyquist rate of $\max |\mathcal{Z}|$. Thus finding a feasible sampling pattern for multicoset sampling is closely connected to solving $\mathfrak{i}(\mathcal{D})$.

While the sampling pattern given in this section is motivated from [14] and [15], follow up works for e.g. in [16], [17] build on the techniques described here.

## B. Fuglede conjecture

Another motivation is related to a conjecture of Fuglede [1]. We say that a set $\mathcal{J} \subseteq \mathbb{Z}_{N}$ tiles $\mathbb{Z}_{N}$ if $\mathcal{J}$ together with its translates forms a disjoint cover of $\mathbb{Z}_{N}$. More precisely, $\mathcal{J}$ tiles $\mathbb{Z}_{N}$ if there exists a set $\mathcal{K} \in \mathbb{Z}_{N}$ (representing the set of translates) such that $\mathcal{J}+\mathcal{K}=\{j+k: j \in \mathcal{J}, k \in \mathcal{K}\}=\mathbb{Z}_{N}$. We can write this as

$$
\begin{equation*}
1_{\mathcal{J}} * 1_{\mathcal{K}}=1_{\mathbb{Z}_{N}}, \text { or } h_{\mathcal{J}} h_{\mathcal{K}}=\delta_{N} \tag{4}
\end{equation*}
$$

where $\delta_{N}$ is the canonical unit vector in $\mathbb{C}^{N}$ with a 1 in the leading position.

Next, $\mathcal{J} \subseteq \mathbb{Z}_{N}$ is called a spectral set if there exists a square unitary submatrix of $\mathcal{F}$ with columns indexed by $\mathcal{J}$ (when we say "unitary" we mean up to scaling).

A conjecture of Fuglede [1], for $\mathbb{Z}_{N}$, is:
Conjecture 1: (Spectral iff Tiling) A set $\mathcal{J} \subseteq \mathbb{Z}_{N}$ is spectral if and only if $\mathcal{J}$ tiles $\mathbb{Z}_{N}$.
See [6], [7], [18]-[22] for some discussion on this conjecture.
Let us do a straightforward analysis of Conjecture 1: starting with a spectral set $\mathcal{J}$, to prove that $\mathcal{J}$ is a tiling set we need to find the set of translates $\mathcal{K}$ such that (4) holds, thus $h_{\mathcal{K}}$ must vanish on $\mathbb{Z}_{N} \backslash\{0\}$ wherever $h_{\mathcal{J}}$ does not. In particular, finding a $\mathcal{K}$ is equivalent to solving the zero set problem; and finding such an idempotent - or the inability to find one would hopefully give insights into the validity of Conjecture 1 , at least in the direction spectral $\Longrightarrow$ tiling.

## C. A case for non prime powers

One particular solution of $\mathfrak{i}(\mathcal{D})$ for $N=p^{M}$ can readily be obtained with elementary means (proof is skipped for brevity). The complete solution space for the prime power case can also
be characterized as in [2]. What is the situation when $N$ has more than one prime factor?

In addition to natural interest, such a generalization is also pertinent if we consider the motivating problems introduced in Section II. Take the problem of sampling two- (or higher) dimensional signals with a fragmented spectrum, where each fragment occupies a single cell of an integer lattice. In the notation of Section II

$$
\begin{aligned}
& \mathcal{F} f\left(s_{1}, s_{2}\right)=0 \\
& \quad \text { when }\left(s_{1}, s_{2}\right) \notin \bigcup_{(m, n) \in \mathfrak{F}}[m, m+1] \times[n, n+1] .
\end{aligned}
$$

Here $\mathfrak{F} \subseteq \mathbb{Z}_{N} \times \mathbb{Z}_{M}$ indicates the locations of the spectral fragments (see Figure 4). An analysis similar to the prime power case leads to two dimensional idempotents $h \in \mathbb{C}^{N \times M}$. When $N$ and $M$ are coprime, the entries in $h$ can be reindexed to an $N M$ - length vector that is an idempotent in $\mathbb{C}^{N M}$. See the Prime Factor algorithm or Good's algorithm ( [23], [24]) for details on using the Chinese remainder theorem for the re-indexing. For such scenarios, the idempotent to be reconstructed is associated with more than one prime factor.


Fig. 4: Example fragmented spectrum of a 2-D signal. The Fourier transform $\mathcal{F} f\left(s_{1}, s_{2}\right)$ is non zero only in the shaded regions.

As for the sampling scenario, a solution of the zero set problem when $N$ has more than one prime factor is relevant to Fuglede's conjecture as well. Fuglede's conjecture is known to be true when $N$ is a prime power, [21], and the challenge is to make progress when $N$ has more than a single prime factor. See [6], [7] for some attempts to generalize the conjecture for arbitrary $N$.

However, unlike the prime power case, there may not be any nontrivial solution to $\mathfrak{i}_{N}(\mathcal{D})$ for an arbitrary $N$. For example, let $N=6$ and $\mathcal{Z}=\{2,3,4\}$. An exhaustive search shows that there is no idempotent $h$ on $\mathbb{Z}_{6}$ that vanishes only on $\mathcal{Z}$ : in fact the only solution to $\mathfrak{i}(\mathcal{D})$ is $\mathbb{Z}_{6}$ (the corresponding idempotent is $\delta$, which vanishes on all non zero indices). This can also be verified from Corollary 2.

## III. EXTENSION TO NON PRIME POWER CASE

We note some simple properties of solutions to $\mathfrak{i}(\mathcal{D})$ that we will use later in the proofs.

Lemma 2: With $\mathfrak{i}(\mathcal{D})$ as defined in the Introduction,

1) The set $\mathcal{J}=\mathbb{Z}_{N}$ is always a solution to $\mathfrak{i}(\mathcal{D})$, for any $\mathcal{D}$.
2) If $\mathcal{J}$ is a solution, so is any translate of $\mathcal{J}, \bmod N$.
3) If $\mathcal{J}_{1}$ and $\mathcal{J}_{2}$ are two disjoint solutions, so is $\mathcal{J}_{1} \cup \mathcal{J}_{2}$.
4) Any solution to $\mathfrak{i}(\mathcal{D})$ is also a solution to $\mathfrak{i}\left(\mathcal{D}^{\prime}\right)$ for any $\mathcal{D}^{\prime} \subseteq \mathcal{D}$.
5) If $\mathcal{J}$ and $\mathcal{J}^{\prime}$ are solutions with $\mathcal{J}^{\prime} \supseteq \mathcal{J}$, then $\mathcal{J}^{\prime} \backslash \mathcal{J}$ is also a solution.
We omit the proof, which uses only the definition and basic properties of the Fourier transform.

Assume, for simplicity, that $N=p q$, so only the two prime factors $p$ and $q$. Suppose that $1 \in \mathcal{D}$, first observe that any solution $h$ to $\mathfrak{i}(\mathcal{D})$ satisfies

$$
\begin{equation*}
h \cdot 1_{\mathscr{A}_{N}(1)}=0 . \tag{5}
\end{equation*}
$$

Recall that $\mathcal{F} h=1_{\mathcal{J}}$, and set $c=\mathcal{F}^{-1} 1_{\mathscr{A}_{N}(1)}$. Note that the entries in $c$ are, by definition of the inverse Fourier transform,

$$
c(k)=\frac{1}{N} \sum_{\substack{n \in \mathbb{Z}_{q} \\(n, N)=1}} \exp (2 \pi i n k / q)
$$

This is Ramanujan's sum; see his original paper [10], for example. As often presented, Ramanujan's sum is

$$
\begin{equation*}
\mathfrak{c}_{N}(k)=\sum_{\substack{n \in \mathbb{Z}_{q} \\(n, N)=1}} \cos (2 \pi n k / N)=\sum_{\substack{n \in \mathbb{Z}_{q} \\(n, N)=1}} \exp (2 \pi i n k / N) \tag{6}
\end{equation*}
$$

where $N$ and $k$ are positive integers. Ramanujan's sums have also recently been used for signal processing, see [11], [12]. For $k \in[0: N-1]$, we can also interpret $\mathfrak{c}_{N}$ as the inverse Fourier transform of the $N$-length (scaled) indicator

$$
\mathbf{1}_{N}(n)= \begin{cases}N, & \text { if }(n, N)=1 \\ 0, & \text { otherwise }\end{cases}
$$

We shall need the following two properties:
(i) When $N=p^{m}$ is a prime power,

$$
\mathfrak{c}_{p^{m}}(k)= \begin{cases}0 & \text { if } p^{m-1} \nmid k,  \tag{7}\\ -p^{m-1} & \text { if } p^{m-1} \mid k \text { and } p^{m} \nmid k \\ \phi\left(p^{m}\right) & \text { if } p^{m} \mid k\end{cases}
$$

where $\phi$ is the Euler totient function [25].
(ii) For any divisor $d$ of $N$,

$$
\sum_{\substack{n \in \mathbb{Z}_{q} \\(n, N)=d}} \exp (2 \pi i n k / N)=\mathfrak{c}_{d^{\prime}}(k)
$$

where $d^{\prime}=N / d$.
We also need the important multiplicative property:

$$
\mathfrak{c}_{p q}(n)=\mathfrak{c}_{p}(n) \mathfrak{c}_{q}(n), \quad \text { if } p, q \text { are co-prime }
$$

Then for $N=p q$ we have

$$
\mathfrak{c}_{N}(j)=\mathfrak{c}_{p q}(j)=\left\{\begin{array}{l}
-(p-1) \text { if } j=p, 2 p, \ldots,(q-1) p  \tag{8}\\
-(q-1) \text { if } j=q, 2 q, \ldots,(p-1) q \\
(p-1)(q-1) \text { if } j=0 \\
1 \text { otherwise }
\end{array}\right.
$$

for $0 \leq j \leq p q-1$.
Going back to (5) and taking Fourier transforms on both sides of $h \cdot 1_{\mathscr{A}_{N}(1)}=0$, we get for any $n \in \mathbb{Z}_{N}$,

$$
\begin{equation*}
\sum_{j \in \mathcal{J}} c(n+j)=0 \quad \text { or } \quad \sum_{j \in \mathcal{K}} c(j)=0 \tag{9}
\end{equation*}
$$

for any set $\mathcal{K}$ that is a translate of $\mathcal{J}$.
We can now move to the main result of this section, characterizing the solutions to $\mathfrak{i}_{p q}(\{1\})$ :

Theorem 1: Suppose $p$ and $q$ are distinct primes, and let $A_{p}=\{0, p, 2 p, \ldots,(q-1) p\}$ and $A_{q}=\{0, q, 2 q, \ldots,(p-$ 1) $q\}$ be subsets of $\mathbb{Z}_{p q}$ containing multiples of $p$ and multiples of $q$, respectively. Then any solution to $\mathfrak{i}_{p q}(\{1\})$ is either a (disjoint) union of translates of $A_{p}$ or a (disjoint) union of translates of $A_{q}$.

Proof: First, we note that both $A_{p}$ and $A_{q}$ are solutions to $\mathfrak{i}(\{1\})$ : we can easily verify from the definition of the discrete Fourier transform

$$
h_{A_{p}}=\mathcal{F}^{-1} \underline{1}_{A_{p}}=\frac{1}{p} \underline{1}_{A_{q}}, \quad \text { and } h_{A_{q}}=\mathcal{F}^{-1} \underline{1}_{A_{q}}=\frac{1}{q} \underline{1}_{A_{p}},
$$

and so both $h_{A_{p}}$ and $h_{A_{q}}$ vanish on indices coprime to $p q$.
Now we prove that any solution to $\mathfrak{i}(\{1\})$ is a disjoint union of translates of either $A_{p}$ or $A_{q}$. Consider any solution $\mathcal{J}$ to $\mathfrak{i}(\{1\})$ in $\mathbb{Z}_{p q}$, and assume, by a suitable translation, that $0 \in \mathcal{J}$. The other elements of $\mathcal{J}$ are either multiples of $p$, multiples of $q$, or coprime to $p q$. Let
$\alpha=\left|\mathcal{J} \cap A_{p} \backslash\{0\}\right|, \quad \beta=\left|\mathcal{J} \cap A_{q} \backslash\{0\}\right|, \quad \gamma=|\mathcal{J}|-\alpha-\beta-1$.
Here $\alpha, \beta$ are the number of nonzero elements of $\mathcal{J}$ that are multiples of $p$ and $q$, respectively, and $\gamma$ is is the number of elements of $\mathcal{J}$ coprime to $p q$. Then we must have, by (8) and (9) (with $n=0$ )

$$
\begin{equation*}
(p-1)(q-1)-\alpha(p-1)-\beta(q-1)+\gamma=0 \tag{10}
\end{equation*}
$$

Note that $\alpha \leq\left|A_{p} \backslash\{0\}\right|=q-1, \beta \leq\left|A_{q} \backslash\{0\}\right|=p-1$, and $\gamma \leq \phi(p q)=(p-1)(q-1)$. We will next argue that either $\alpha=q-1$ or $\beta=p-1$ (i.e. either $\alpha$ or $\beta$ take their largest possible values); thus establishing that $\mathcal{J}$ contains either $A_{p}$ or $A_{q}$.

If $\alpha=q-1$, then $A_{p} \subseteq \mathcal{J}$ and we are done. So suppose that $\alpha<q-1$. Then there exists a non zero multiple of $p$, say $a p$, such that $a p \notin \mathcal{J}$. Now consider the set $\mathcal{J}^{\prime}=\tau^{a p} \mathcal{J}$ obtained by translating $\mathcal{J}$ by $a p$. We first make these simple observations:

1) By construction, $0 \notin \mathcal{J}^{\prime}$. The element $0 \in \mathcal{J}$ translated by $a p$ results in a non zero multiple of $p$.
2) Translating non zero multiples of $p$ in $\mathcal{J}$ by $a p$ results in indices that are non zero multiples of $p$.
3) Translating non zero multiples of $q$ in $\mathcal{J}$ by $a p$ results in indices that are co-prime to $p q$.
4) Translating indices in $\mathcal{J}$ that are co-prime to $p q$ by $a p$ could potentially result in indices that are multiples of $q$. Let $\gamma_{1} \leq \gamma$ be the number of such indices.
Applying (9) with $n=-a p$, and using (8) results in

$$
\begin{equation*}
-(\alpha+1)(p-1)-\gamma_{1}(q-1)+\beta+\left(\gamma-\gamma_{1}\right)=0 \tag{11}
\end{equation*}
$$

Subtracting (11) from (10) results in

$$
\left(p-1-\beta+\gamma_{1}\right) q=0 \quad \text { or } p-1+\gamma_{1}=\beta
$$

Since $\beta \leq p-1$, the above equality only possible when $\beta=$ $p-1$ and $\gamma_{1}=0$. In particular $\beta=p-1$ implies that $A_{q} \subseteq \mathcal{J}$. Thus either $A_{p} \subseteq \mathcal{J}$ or $A_{q} \subseteq \mathcal{J}$.

Now we argue as follows: isolate either $A_{p}$ or $A_{q}$ from $\mathcal{J}$, as in

$$
\mathcal{J}_{1}=\left\{\begin{array}{l}
A_{p} \text { if } A_{p} \subseteq \mathcal{J} \\
A_{q} \text { otherwise }
\end{array}\right.
$$

We note, as before, that $\mathcal{J} \backslash \mathcal{J}_{1}$ is a solution to $\mathfrak{i}(\{1\})$. Applying the argument repeatedly, we can write

$$
\mathcal{J}=\mathcal{J}_{1} \cup \mathcal{J}_{2} \cup \mathcal{J}_{3} \ldots
$$

where each $\mathcal{J}_{i}$ is a translate of $A_{p}$ or $A_{q}$, and all the $\mathcal{J}_{i}$ are disjoint. We will argue next that all of the $\mathcal{J}_{i}$ are translates of the same set: either $A_{p}$ or $A_{q}$. Suppose that $\mathcal{J}_{1}=\{0, p, 2 p, \ldots,(q-1) p\}$ and $\mathcal{J}_{k}=\{l, l+q, l+2 q, \ldots, l+$ $(p-1) q\}$ is a translate of $A_{q}$. Then there exists an element of $\mathcal{J}_{k}$ that is a multiple of $p$ : this is obtained by solving the congruence

$$
l+x q \equiv 0 \quad \bmod p
$$

in $x$. Thus $\mathcal{J}_{1}$ and $\mathcal{J}_{k}$ have a common element, contradicting their disjointedness.

So $\mathcal{J}$ is a disjoint union of translates of $A_{p}$, or a disjoint union of translates of $A_{q}$, proving the theorem.

Corollary 1: If $p$ and $q$ are distinct primes, any solution to $\mathfrak{i}_{p q}(\{1\})$ is a solution to either $\mathfrak{i}_{p q}(\{1, p\})$ or $\mathfrak{i}_{p q}(\{1, q\})$. In particular, there is no idempotent that vanishes only on $\mathscr{A}_{p q}(1)$. Proof: Recall that

$$
h_{A_{p}}=\mathcal{F}^{-1} \underline{1}_{A_{p}}=\frac{1}{p} \underline{1}_{A_{q}}, \quad \text { and } h_{A_{q}}=\mathcal{F}^{-1} \underline{1}_{A_{q}}=\frac{1}{q} \underline{1}_{A_{p}}
$$

Thus $h_{A_{p}}$ is nonzero only on multiples of $q$, i.e. it vanishes on $\mathscr{A}(1) \cup \mathscr{A}(p)$, and so $A_{p}$ is a solution to $\mathfrak{i}(\{1, p\})$ and (similarly) $A_{q}$ is a solution to $\mathfrak{i}(\{1, q\})$. From Theorem 1, any solution to $\mathfrak{i}(\{1\})$ is a disjoint union of translates of $A_{p}$ or $A_{q}$. From this and Lemma 2 properties 2) and 3), it follows that any solution to $\mathfrak{i}(\{1\})$ is a solution to either $\mathfrak{i}(\{1, p\})$ or $\mathfrak{i}(\{1, q\})$.

Corollary 2: If $p$ and $q$ are distinct primes, then the only solution to $\mathfrak{i}_{p q}(\{p, q\})$ is $\mathbb{Z}_{N}$.

Proof: Let $\mathcal{J}$ be a solution to $\mathfrak{i}_{p q}(\{p, q\})$, and $\mathcal{K}$ any solution to $\mathfrak{i}_{p q}(\{1\})$. Then $h_{\mathcal{J}} \cdot h_{\mathcal{K}}$ must vanish at all indices in $\mathbb{Z} \backslash\{0\}$. Note that $h_{\mathcal{J}}(0)=|\mathcal{J}| / N$ and $h_{\mathcal{K}}(0)=|\mathcal{K}| / N$ (here $N=p q$ ). Then we must have $h_{\mathcal{J}} \cdot h_{\mathcal{K}}=|\mathcal{J}| \mathcal{K} \mid \delta_{p q} / N^{2}$. Using that $\mathcal{F} h_{\mathcal{J}}=1_{\mathcal{J}}$ and $\mathcal{F} h_{\mathcal{K}}=1_{\mathcal{K}}$, we obtain

$$
1_{\mathcal{J}} * 1_{\mathcal{K}}=N \mathcal{F}\left(h_{\mathcal{J}} \cdot h_{\mathcal{K}}\right)=\frac{|\mathcal{J}||\mathcal{K}|}{N} 1_{\mathbb{Z}_{N}}
$$

In particular, since $1_{\mathcal{J}}, 1_{\mathcal{K}}$ are indicators, the convolution $1_{\mathcal{J}} *$ $1_{\mathcal{K}}$ can only have integer entries; and so we must have that ${ }_{65}^{N}=p q$ divides $|\mathcal{J}||\mathcal{K}|$. This must be true for any $\mathcal{K}$ that
is solution to $\mathfrak{i}_{p q}(\{1\})$. In particular, we can use $\mathcal{K}=A_{p}$ or $\mathcal{K}=A_{q}$ to conclude that $|\mathcal{J}|=p q$, or in other words $\mathcal{J}=\mathbb{Z}_{N}$.

## IV. A bound on the smallest solution to $\mathfrak{i}(\mathcal{D})$

Recall from our sampling based motivation for the zero set problem in Section II that the average sampling rate of the multi-coset scheme discussed is $|\mathcal{J}|$, which we would like to keep as small as possible. Naturally, instead of asking for all possible idempotents with a given zero set (as in $\mathfrak{i}(\mathcal{D})$ ) we can ask for the smallest solution to $\mathfrak{i}(\mathcal{D})$. This leads to the definition

$$
\Xi(\mathcal{D})=\arg \min \{|\mathcal{J}|: \mathcal{J} \text { solves } \mathfrak{i}(\mathcal{D})\}
$$

This quantity, in some sense, characterizes the smallest possible sampling rate achievable by the multicoset scheme.

In this section, we derive some simple bounds on $\Xi(\mathcal{D})$, and look at some connections to difference graphs defined in an earlier work [6]. To define the bound, given a set of divisors $\mathcal{D}$, recall that the corresponding zero set $\mathcal{Z}$ is defined as in (2) by including all indices whose gcd with $N$ is in $\mathcal{D}$. In the (inverse) DFT matrix, suppose we remove all rows except those with indices in $\mathcal{Z} \cup\{0\}$, to obtain the matrix $\mathcal{F}_{\mathcal{Z} \cup\{0\}}^{-1}$.

Note that if $\mathcal{J}$ is a solution to $\mathfrak{i}(\mathcal{D})$, then $\mathcal{F}^{-1} 1_{\mathcal{J}}$ vanishes on $\mathcal{Z}$ by definition. In addition, we also have $\mathcal{F}^{-1} 1_{\mathcal{J}}(0)=$ $|\mathcal{J}| / N \neq 0$, so that we may say

$$
\mathcal{F}_{\mathcal{Z} \cup\{0\}}^{-1}\left(1_{\mathcal{J}} N /|\mathcal{J}|\right)=\delta,
$$

where, as usual, $\delta$ is standard unit vector with 1 in the topmost position. Now we construct a lower-bound to $\Xi(\mathcal{D})$ by replacing $1_{\mathcal{J}} N /|\mathcal{J}|$ with $x$ in the above equation:

$$
\xi(\mathcal{D})=\min _{x \in \mathbb{C}^{N}}\left\{\|x\|_{0}: \mathcal{F}_{\mathcal{Z} \cup\{0\}}^{-1} x=\delta\right\}
$$

where $\|x\|_{0}$ is the number of non zero entries in $x$. The above formulation asks us to find the sparsest solution to a system of linear equations, for which we can potentially apply standard algorithms like Orthogonal Matching Pursuit or Basis Pursuit [26]-[28]. By construction, we have

$$
\begin{equation*}
\xi(\mathcal{D}) \leq \Xi(\mathcal{D}) \tag{12}
\end{equation*}
$$

## A. Difference graphs

Given a divisor set $\mathcal{D}$ and the corresponding zero set $\mathcal{Z}$ we can define a graph $\mathscr{G}(\mathcal{D})$ with vertex set $\mathbb{Z}_{N}$ and edge between $i, j$ if $(i-j, N) \in \mathcal{Z}$. Such graphs were investigated in our prior work [6], in the context of sampling discrete signals. We explore some connections of the bound $\xi(\mathcal{D})$ to certain graph invariants.

We say that a matrix $M$ fits a graph $G$ if the diagonal entries of $M$ are 1 , and $i j$ entry $M_{i j}$ is zero if $i, j$ are not adjacent in $G$. Recall that the minrank of a graph (over $\mathbb{C}$ )) is the smallest possible rank among all complex matrices that fit $G$ [13], [29], [30].

Start with a divisor set $\mathcal{D}$ and the corresponding zero set $\mathcal{Z}$. Let $x \in \mathbb{C}^{N}$ satisfy $\mathcal{F}_{\mathcal{Z} \cup\{0\}}^{-1} x=\delta$. Construct an $N \times N$
circulant matrix [31] $M$ with first column $\mathcal{F}^{-1} x$. Note that the eigen values of $M$ are the entries of $x$, and consequently the rank of $M$ is $\|x\|_{0}$. Also note that the diagonal entries of $M$ are 1 , and the $i j$ entry is zero if $i-j \in \mathcal{Z}$. Thus $M$ fits $\mathscr{G}^{c}(\mathcal{D})$, and so

$$
\operatorname{minrank}\left(\mathscr{G}^{c}(\mathcal{D})\right) \leq \xi(\mathcal{D})
$$

The bound $\xi(\mathcal{D})$ is similar to minimum circulant rank defined in [32]. We can infact prove that the bound in (12) is tight, suprisingly in the case when $\mathcal{J}$ is spectral.

Lemma 3: Suppose $\mathcal{J}$ is spectral, and $h=\mathcal{F}^{-1} 1_{\mathcal{J}}$, as before. Then the bound in (12) is tight, i.e. $\xi(\mathcal{D}(h))=\Xi(\mathcal{D}(h))$.

Proof: From [6, Lemma 3], if $\mathcal{J}$ is spectral then there exists an independent set in $\mathscr{G}^{c}(\mathcal{D})$ of size $|\mathcal{J}|$. Since minrank is an upper bound on the independence number [13], it follows that $|\mathcal{J}| \leq \operatorname{minrank}\left(\mathscr{G}^{c}(\mathcal{D})\right)$. Combining with (12) we have that

$$
|\mathcal{J}| \leq \operatorname{minrank}\left(\mathscr{G}^{c}(\mathcal{D})\right) \leq \xi(\mathcal{D}) \leq \Xi(\mathcal{D}) \leq|\mathcal{J}|
$$

and so all the inequalities involved are tight.
Another lower bound can be defined using Linear programs [7, Section IV.B]. Note that this bound is non integral, unlike the bound $\xi$ given in this section.

## V. Conclusion

We introduced the zero set problem for convolution idempotents, briefly reviewed the motivations, and presented some results when the ambient dimension is a product of two primes. The connection to results on vanishing sums of roots of unity [9] and tiling [21] need to be explored further. We also gave some bounds on the smallest solution to the zero set problem and explored its connection to minrank of certain graphs defined on $\mathbb{Z}_{N}$. Of interest is to investigate generalizations to arbitrary $N$. Of particular interest is to understand the solution space of $\mathfrak{i}(\mathcal{D})$ when $\mathcal{D}$ corresponds to spectral or tiling sets [21], the conditions under which the bound $\Xi(\mathcal{D})$ can be efficiently computed, and provable algorithms to solve (or approximate) the solutions to $\mathfrak{i}(\mathcal{D})$.

## References

[1] B. Fuglede, "Commuting self-adjoint partial differential operators and a group theoretic problem," Journal of Functional Analysis, vol. 16, no. 1, pp. 101-121, 1974.
[2] A. Siripuram and B. Osgood, "Convolution idempotents with a given zero-set," 2020.
[3] A. V. Oppenheim, Discrete-time signal processing. Pearson Education India, 1999.
[4] B. Osgood, Lectures on the Fourier Transform and Its Applications. American Mathematical Society, 2018.
[5] W. Wu, "Discrete sampling: Generalizations of the Nyquist-Shannon sampling theorem," Ph.D. dissertation, Stanford University, 2010.
[6] A. Siripuram, W. Wu, and B. Osgood, "Discrete sampling: A graph theoretic approach to orthogonal interpolation," IEEE Transactions on Information Theory, 2019.
[7] A. Siripuram and B. Osgood, "Lp relaxations and fuglede's conjecture," in 2018 IEEE International Symposium on Information Theory (ISIT). IEEE, 2018, pp. 2525-2529.
[8] R.-D. Malikiosis and M. N. Kolountzakis, "Fuglede's conjecture on cyclic groups of order $p^{n} q$," arXiv preprint arXiv:1612.01328, 2016.
[9] T. Lam and K. Leung, "On vanishing sums of roots of unity," Journal of algebra, vol. 224, no. 1, pp. 91-109, 2000.
[10] S. Ramanujan, "On certain trigonometrical sums and their applications in the theory of numbers," Trans. Cambridge Philos. Soc, vol. 22, no. 13, pp. 259-276, 1918.
[11] P. Vaidyanathan, "Ramanujan sums in the context of signal process-ing-part i: Fundamentals," IEEE transactions on signal processing, vol. 62, no. 16, pp. 4145-4157, 2014.
[12] _-, "Ramanujan sums in the context of signal processing-part ii: FIR representations and applications," IEEE Transactions on Signal Processing, vol. 62, no. 16, pp. 4158-4172, 2014.
[13] W. Haemers et al., "An upper bound for the shannon capacity of a graph," in Colloq. Math. Soc. János Bolyai, vol. 25, 1978, pp. 267-272.
[14] A. Kohlenberg, "Exact interpolation of band-limited functions," Journal of Applied Physics, vol. 24, no. 12, pp. 1432-1436, 1953.
[15] Y.-P. Lin and P. Vaidyanathan, "Periodically nonuniform sampling of bandpass signals," IEEE Transactions on Circuits and Systems II: Analog and Digital Signal Processing, vol. 45, no. 3, pp. 340-351, 1998.
[16] C. Herley and P. W. Wong, "Minimum rate sampling and reconstruction of signals with arbitrary frequency support," IEEE Transactions on Information Theory, vol. 45, no. 5, pp. 1555-1564, 1999.
[17] R. Venkataramani and Y. Bresler, "Perfect reconstruction formulas and bounds on aliasing error in sub-nyquist nonuniform sampling of multiband signals," IEEE Transactions on Information Theory, vol. 46, no. 6, pp. 2173-2183, 2000.
[18] A. Iosevich, N. Katz, and T. Tao, "The Fuglede spectral conjecture holds for convex planar domains," Mathematical Research Letters, vol. 10, no. 5, pp. 559-569, 2003.
[19] I. Laba, "The spectral set conjecture and multiplicative properties of roots of polynomials," Journal of the London Mathematical Society, vol. 65, no. 03, pp. 661-671, 2002.
[20] T. Tao, "Fuglede's conjecture is false in 5 and higher dimensions," arXiv preprint math/0306134, 2003.
[21] E. M. Coven and A. Meyerowitz, "Tiling the integers with translates of one finite set," Journal of Algebra, vol. 212, no. 1, pp. 161-174, 1999.
[22] D. E. Dutkay and C.-K. LAI, "Some reductions of the spectral set conjecture to integers," in Mathematical Proceedings of the Cambridge Philosophical Society, vol. 156, no. 01. Cambridge Univ Press, 2014, pp. 123-135.
[23] I. J. Good, "The interaction algorithm and practical fourier analysis," Journal of the Royal Statistical Society. Series B (Methodological), pp. 361-372, 1958.
[24] P. Duhamel and M. Vetterli, "Fast fourier transforms: a tutorial review and a state of the art," Signal processing, vol. 19, no. 4, pp. 259-299, 1990.
[25] R. L. Graham, D. E. Knuth, O. Patashnik, and S. Liu, "Concrete mathematics: a foundation for computer science," Computers in Physics, vol. 3, no. 5, pp. 106-107, 1989.
[26] T. T. Cai and L. Wang, "Orthogonal matching pursuit for sparse signal recovery with noise." Institute of Electrical and Electronics Engineers, 2011.
[27] M. Elad, Sparse and redundant representations: from theory to applications in signal and image processing. Springer Science \& Business Media, 2010.
[28] E. J. Candès and M. B. Wakin, "An introduction to compressive sampling [a sensing/sampling paradigm that goes against the common knowledge in data acquisition]," IEEE signal processing magazine, vol. 25, no. 2, pp. 21-30, 2008.
[29] A. Blasiak, R. Kleinberg, and E. Lubetzky, "Broadcasting with side information: Bounding and approximating the broadcast rate," IEEE Transactions on Information Theory, vol. 59, no. 9, pp. 5811-5823, 2013.
[30] B. Bukh and C. Cox, "On a fractional version of haemers' bound," IEEE Transactions on Information Theory, vol. 65, no. 6, pp. 33403348, 2018.
[31] P. J. Davis, Circulant matrices. American Mathematical Soc., 2013.
[32] L. Deaett and S. A. Meyer, "The minimum rank problem for circulants," Linear Algebra and its Applications, vol. 491, pp. 386-418, 2016.

