# Some results on convolution idempotents

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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 hvanishes. 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  $\mathcal{J} \subseteq \mathcal{L}$  vanishing sums of roots of unity [9], though the te employ are different. The proof of Theorem 1 use of Ramanujan sums [10], which were recently use processing [11], [12]. In section IV we invest bounds on the *smallest* solution to  $i(\mathcal{D})$ , and rel minrank [13] of certain graphs introduced in [6]. 1462

omitting N), let (i, N) denote the greatest common divisor (gcd) of i and N, and let

$$\mathscr{A}_N(k) = \{ i \in \mathbb{Z}_N \colon (i, N) = k \}.$$

$$(1)$$

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

$$\mathcal{Z}(h) = \{i \in \mathbb{Z}_N \colon (i, N) \in \mathcal{D}(h)\} = \bigcup_{k \in \mathcal{D}(h)} \mathscr{A}_N(k)$$

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  $\mathcal{Z}(h_{\mathcal{J}}) = \{1, 3, 4, 5, 7\} = \{1, 3, 5, 7\} \cup \{4\} = \mathscr{A}_8(1) \cup \mathscr{A}_8(4)$ ; here  $\mathcal{D}(h_{\mathcal{J}}) = \{1, 4\}.$ 

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  $i_N(\mathcal{D})$ : Given a positive integer N and a set of divisors  $\mathcal{D} \subseteq \mathcal{D}_N$  let

$$\mathcal{Z} = \{ i \in \mathbb{Z}_N \colon (i, N) \in \mathcal{D} \}$$
(2)

Find all index sets  $\mathcal{J}$  such that the idempotent  $h_{\mathcal{J}} = \mathcal{F}^{-1} \mathbf{1}_{\mathcal{J}}$  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  $i(\mathcal{D})$  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 = pq and  $D = \{1\}, \{p, q\}, \{1, p\}$  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  $i(\mathcal{D})$ , and relate it to the **ISIT 2020** 

## 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  $\bigcup_{n \in \mathfrak{F}} [n, n+1]$ . 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

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

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

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.,

$$\mathcal{Z} = \{k_1 - k_2 : k_1, k_2 \in \mathfrak{F}\},\$$
$$\mathcal{D} = \{d : (k, N) = d \text{ for some } k \in \mathfrak{F}.\}$$

We can show ([2, Proposition 2]) that any  $\mathcal{J}$  providing a solution to the zero set problem  $i(\mathcal{D})$  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  $i(\mathcal{D})$ .

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

$$1_{\mathcal{J}} * 1_{\mathcal{K}} = 1_{\mathbb{Z}_N}, \text{ or } h_{\mathcal{J}} h_{\mathcal{K}} = \delta_N,$$
 (4)

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  $\mathbb{Z}_N \setminus \{0\}$  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  $\implies$  tiling.

## C. A case for non prime powers

One particular solution of  $i(\mathcal{D})$  for  $N = p^M$  can readily be obtained with elementary means (proof is skipped for brevity). fragments are intact: aliasing primary occurs in the islands The complete solution space for the prime power case can also 1463 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

$$\mathcal{F}f(s_1, s_2) = 0$$
  
when  $(s_1, s_2) \notin \bigcup_{(m,n) \in \mathfrak{F}} [m, m+1] \times [n, n+1].$ 

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 NM- length vector that is an idempotent in  $\mathbb{C}^{NM}$ . 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  $\mathcal{F}f(s_1, s_2)$  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  $i_N(\mathcal{D})$  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  $i(\mathcal{D})$  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 i(D) that we will use later in the proofs.

*Lemma 2:* With i(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  $\mathcal{J}$ , mod N.
- 3) If  $\mathcal{J}_1$  and  $\mathcal{J}_2$  are two disjoint solutions, so is  $\mathcal{J}_1 \cup \mathcal{J}_2$ .
- Any solution to i(D) is also a solution to i(D') for any D' ⊆ D.
- 5) If  $\mathcal{J}$  and  $\mathcal{J}'$  are solutions with  $\mathcal{J}' \supseteq \mathcal{J}$ , then  $\mathcal{J}' \setminus \mathcal{J}$  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 = pq, 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

$$h \cdot 1_{\mathscr{A}_N(1)} = 0. \tag{5}$$

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,

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

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

$$\mathfrak{c}_N(k) = \sum_{\substack{n \in \mathbb{Z}_q \\ (n,N)=1}} \cos(2\pi nk/N) = \sum_{\substack{n \in \mathbb{Z}_q \\ (n,N)=1}} \exp(2\pi ink/N),$$
(6)

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  $c_N$  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,

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

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

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

where d' = N/d.

We also need the important multiplicative property:

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

Then for N = pq we have

$$\mathfrak{c}_{N}(j) = \mathfrak{c}_{pq}(j) = \begin{cases} -(p-1) \text{ if } j = p, 2p, \dots, (q-1)p \\ -(q-1) \text{ if } j = q, 2q, \dots, (p-1)q \\ (p-1)(q-1) \text{ if } j = 0 \\ 1 \text{ otherwise} \end{cases},$$
(8)

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for  $0 \le j \le pq - 1$ .

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

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

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  $i_{pq}(\{1\})$ :

Theorem 1: Suppose p and q are distinct primes, and let  $A_p = \{0, p, 2p, \ldots, (q-1)p\}$  and  $A_q = \{0, q, 2q, \ldots, (p-1)q\}$  be subsets of  $\mathbb{Z}_{pq}$  containing multiples of p and multiples of q, respectively. Then any solution to  $i_{pq}(\{1\})$  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  $i(\{1\})$ : we can easily verify from the definition of the discrete Fourier transform

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

and so both  $h_{A_p}$  and  $h_{A_q}$  vanish on indices coprime to pq.

Now we prove that any solution to  $i(\{1\})$  is a disjoint union of translates of either  $A_p$  or  $A_q$ . Consider any solution  $\mathcal{J}$ to  $i(\{1\})$  in  $\mathbb{Z}_{pq}$ , 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 pq. Let

$$\alpha = |\mathcal{J} \cap A_p \setminus \{0\}|, \quad \beta = |\mathcal{J} \cap A_q \setminus \{0\}|, \quad \gamma = |\mathcal{J}| - \alpha - \beta - 1.$$

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 pq. Then we must have, by (8) and (9) (with n = 0)

$$(p-1)(q-1) - \alpha(p-1) - \beta(q-1) + \gamma = 0.$$
 (10)

Note that  $\alpha \leq |A_p \setminus \{0\}| = q-1$ ,  $\beta \leq |A_q \setminus \{0\}| = p-1$ , and  $\gamma \leq \phi(pq) = (p-1)(q-1)$ . 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 ap, such that  $ap \notin \mathcal{J}$ . Now consider the set  $\mathcal{J}' = \tau^{ap} \mathcal{J}$ obtained by translating  $\mathcal{J}$  by ap. We first make these simple observations:

- 1) By construction,  $0 \notin \mathcal{J}'$ . The element  $0 \in \mathcal{J}$  translated by *ap* results in a non zero multiple of *p*.
- 2) Translating non zero multiples of p in  $\mathcal{J}$  by ap results in indices that are non zero multiples of p.
- 3) Translating non zero multiples of q in  $\mathcal{J}$  by ap results in indices that are co-prime to pq.
- Translating indices in *J* that are co-prime to pq by ap could potentially result in indices that are multiples of q. Let γ<sub>1</sub> ≤ γ be the number of such indices.

Applying (9) with n = -ap, and using (8) results in

$$-(\alpha + 1)(p - 1) - \gamma_1(q - 1) + \beta + (\gamma - \gamma_1) = 0.$$

Subtracting (11) from (10) results in

$$(p - 1 - \beta + \gamma_1)q = 0$$
 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

$$\mathcal{J}_1 = \begin{cases} A_p \text{ if } A_p \subseteq \mathcal{J} \\ A_q \text{ otherwise} \end{cases}$$

We note, as before, that  $\mathcal{J} \setminus \mathcal{J}_1$  is a solution to  $i(\{1\})$ . Applying the argument repeatedly, we can write

$$\mathcal{J} = \mathcal{J}_1 \cup \mathcal{J}_2 \cup \mathcal{J}_3 \dots$$

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  $\mathcal{J}_1 = \{0, p, 2p, \dots, (q-1)p\}$  and  $\mathcal{J}_k = \{l, l+q, l+2q, \dots, l+(p-1)q\}$  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

$$l + xq \equiv 0 \mod p$$

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  $i_{pq}(\{1\})$  is a solution to either  $i_{pq}(\{1, p\})$  or  $i_{pq}(\{1, q\})$ . In particular, there is no idempotent that vanishes only on  $\mathscr{A}_{pq}(1)$ .

*Proof:* Recall that

$$h_{A_p} = \mathcal{F}^{-1} \underline{1}_{A_p} = \frac{1}{p} \underline{1}_{A_q}, \text{ and } h_{A_q} = \mathcal{F}^{-1} \underline{1}_{A_q} = \frac{1}{q} \underline{1}_{A_p}.$$

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  $i(\{1, p\})$  and (similarly)  $A_q$  is a solution to  $i(\{1, q\})$ . From Theorem 1, any solution to  $i(\{1\})$  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  $i(\{1\})$  is a solution to either  $i(\{1, p\})$  or  $i(\{1, q\})$ .

Corollary 2: If p and q are distinct primes, then the only solution to  $i_{pq}(\{p,q\})$  is  $\mathbb{Z}_N$ .

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

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

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 (11) N = pq divides  $|\mathcal{J}||\mathcal{K}|$ . This must be true for any  $\mathcal{K}$  that is solution to  $i_{pq}(\{1\})$ . In particular, we can use  $\mathcal{K} = A_p$  or  $\mathcal{K} = A_q$  to conclude that  $|\mathcal{J}| = pq$ , 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  $i(\mathcal{D})$ ) we can ask for the smallest solution to  $i(\mathcal{D})$ . 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  $\mathcal{F}^{-1}\mathbf{1}_{\mathcal{J}}$  vanishes on  $\mathcal{Z}$  by definition. In addition, we also have  $\mathcal{F}^{-1}\mathbf{1}_{\mathcal{J}}(0) = |\mathcal{J}|/N \neq 0$ , 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} \{ \|x\|_0 : \mathcal{F}_{\mathcal{Z} \cup \{0\}}^{-1} x = \delta \},\$$

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

$$\xi(\mathcal{D}) \le \Xi(\mathcal{D}). \tag{12}$$

## 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 ij entry  $M_{ij}$  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  $N \times N$ 1466

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 ij entry is zero if  $i - j \in \mathbb{Z}$ . Thus M fits  $\mathscr{G}^c(\mathcal{D})$ , and so

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

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  $|\mathcal{J}| \leq \text{minrank}(\mathscr{G}^{c}(\mathcal{D}))$ . Combining with (12) we have that

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

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  $i(\mathcal{D})$  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  $i(\mathcal{D})$ .

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