



# Optimizing Sparsity over Lattices and Semigroups

Iskander Aliev<sup>1</sup>, Gennadiy Averkov<sup>2(✉)</sup>, Jesús A. De Loera<sup>3</sup>, and Timm Oertel<sup>1</sup>

<sup>1</sup> Cardiff University, Cardiff, UK

<sup>2</sup> Brandenburg University of Technology Cottbus-Senftenberg, Senftenberg, Germany  
averkov@b-tu.de

<sup>3</sup> University of California, Davis, USA

**Abstract.** Motivated by problems in optimization we study the *sparsity* of the solutions to systems of linear Diophantine equations and linear integer programs, i.e., the number of non-zero entries of a solution, which is often referred to as the  $\ell_0$ -norm. Our main results are improved bounds on the  $\ell_0$ -norm of sparse solutions to systems  $A\mathbf{x} = \mathbf{b}$ , where  $A \in \mathbb{Z}^{m \times n}$ ,  $\mathbf{b} \in \mathbb{Z}^m$  and  $\mathbf{x}$  is either a general integer vector (lattice case) or a non-negative integer vector (semigroup case). In the lattice case and certain scenarios of the semigroup case, we give polynomial time algorithms for computing solutions with  $\ell_0$ -norm satisfying the obtained bounds.

## 1 Introduction

This paper discusses the problem of finding sparse solutions to systems of linear Diophantine equations and integer linear programs. We investigate the  $\ell_0$ -norm  $\|\mathbf{x}\|_0 := |\{i : x_i \neq 0\}|$ , a function widely used in the theory of *compressed sensing* [6, 9], which measures the sparsity of a given vector  $\mathbf{x} = (x_1, \dots, x_n)^\top \in \mathbb{R}^n$  (it is clear that the  $\ell_0$ -norm is actually not a norm).

Sparsity is a topic of interest in several areas of optimization. The  $\ell_0$ -norm minimization problem over reals is central in the theory of the classical compressed sensing, where a linear programming relaxation provides a guaranteed approximation [8, 9]. Support minimization for solutions to Diophantine equations is relevant for the theory of compressed sensing for discrete-valued signals [11, 12, 17]. There is still little understanding of discrete signals in the compressed sensing paradigm, despite the fact that there are many applications in which the signal is known to have discrete-valued entries, for instance, in wireless communication [22] and the theory of error-correcting codes [7]. Sparsity was also investigated in integer optimization [1, 10, 20], where many combinatorial optimization problems have useful interpretations as sparse semigroup problems. For example, the edge-coloring problem can be seen as a problem in the semigroup generated by matchings of the graph [18]. Our results provide natural out-of-the-box sparsity bounds for problems with linear constraints and integer variables in a general form.

## 1.1 Lattices: Sparse Solutions of Linear Diophantine Systems

Each integer matrix  $A \in \mathbb{Z}^{m \times n}$  determines the lattice  $\mathcal{L}(A) := \{Ax : x \in \mathbb{Z}^n\}$  generated by the columns of  $A$ . By an easy reduction via row transformations, we may assume without loss of generality that the rank of  $A$  is  $m$ .

Let  $[n] := \{1, \dots, n\}$  and let  $\binom{[n]}{m}$  be the set of all  $m$ -element subsets of  $[n]$ . For  $\gamma \subseteq [n]$ , consider the  $m \times |\gamma|$  submatrix  $A_\gamma$  of  $A$  with columns indexed by  $\gamma$ . One can easily prove that the determinant of  $\mathcal{L}(A)$  is equal to

$$\gcd(A) := \gcd \left\{ \det(A_\gamma) : \gamma \in \binom{[n]}{m} \right\}.$$

Since  $\mathcal{L}(A_\gamma)$  is the lattice spanned by the columns of  $A$  indexed by  $\gamma$ , it is a sublattice of  $\mathcal{L}(A)$ . We first deal with a natural question: *Can the description of a given lattice  $\mathcal{L}(A)$  in terms of  $A$  be made sparser by passing from  $A$  to  $A_\gamma$  with  $\gamma$  having a smaller cardinality than  $n$  and satisfying  $\mathcal{L}(A) = \mathcal{L}(A_\gamma)$ ?* That is, we want to discard some of the columns of  $A$  and generate  $\mathcal{L}(A)$  by  $|\gamma|$  columns with  $|\gamma|$  being possibly small.

For stating our results, we need several number-theoretic functions. Given  $z \in \mathbb{Z}_{>0}$ , consider the prime factorization  $z = p_1^{s_1} \cdots p_k^{s_k}$  with pairwise distinct prime factors  $p_1, \dots, p_k$  and their multiplicities  $s_1, \dots, s_k \in \mathbb{Z}_{>0}$ . Then the number of prime factors  $\sum_{i=1}^k s_i$  counting the multiplicities is denoted by  $\Omega(z)$ . Furthermore, we introduce  $\Omega_m(z) := \sum_{i=1}^k \min\{s_i, m\}$ . That is, by introducing  $m$  we set a threshold to account for multiplicities. In the case  $m = 1$  we thus have  $\omega(z) := \Omega_1(z) = k$ , which is the number of prime factors in  $z$ , not taking the multiplicities into account. The functions  $\Omega$  and  $\omega$  are called *prime  $\Omega$ -function* and *prime  $\omega$ -function*, respectively, in number theory [15]. We call  $\Omega_m$  the *truncated prime  $\Omega$ -function*.

**Theorem 1.** *Let  $A \in \mathbb{Z}^{m \times n}$ , with  $m \leq n$ , and let  $\tau \in \binom{[n]}{m}$  be such that the matrix  $A_\tau$  is non-singular. Then the equality  $\mathcal{L}(A) = \mathcal{L}(A_\tau)$  holds for some  $\gamma$  satisfying  $\tau \subseteq \gamma \subseteq [n]$  and*

$$|\gamma| \leq m + \Omega_m \left( \frac{|\det(A_\tau)|}{\gcd(A)} \right). \quad (1)$$

*Given  $A$  and  $\tau$ , the set  $\gamma$  can be computed in polynomial time.*

One can easily see that  $\omega(z) \leq \Omega_m(z) \leq \Omega(z) \leq \log_2(z)$  for every  $z \in \mathbb{Z}_{>0}$ . The estimate using  $\log_2(z)$  gives a first impression on the quality of the bound (1). It turns out, however, that  $\Omega_m(z)$  is much smaller on the average. Results in number theory [15, §22.10] show that the average values  $\frac{1}{z}(\omega(1) + \cdots + \omega(z))$  and  $\frac{1}{z}(\Omega(1) + \cdots + \Omega(z))$  are of order  $\log \log z$ , as  $z \rightarrow \infty$ .

As an immediate consequence of Theorem 1 we obtain

**Corollary 2.** *Consider the linear Diophantine system*

$$Ax = b, \quad x \in \mathbb{Z}^n \quad (2)$$

with  $A \in \mathbb{Z}^{m \times n}$ ,  $\mathbf{b} \in \mathbb{Z}^m$  and  $m \leq n$ . Let  $\tau \in \binom{[n]}{m}$  be such that the  $m \times m$  matrix  $A_\tau$  is non-singular. If (2) is feasible, then (2) has a solution  $\mathbf{x}$  satisfying the sparsity bound

$$\|\mathbf{x}\|_0 \leq m + \Omega_m \left( \frac{|\det(A_\tau)|}{\gcd(A)} \right).$$

Under the above assumptions, for given  $A, \mathbf{b}$  and  $\tau$ , such a sparse solution can be computed in polynomial time.

From the optimization perspective, Corollary 2 deals with the problem

$$\min \{ \|\mathbf{x}\|_0 : A\mathbf{x} = \mathbf{b}, \mathbf{x} \in \mathbb{Z}^n \}$$

of minimization of the  $\ell_0$ -norm over the affine lattice  $\{\mathbf{x} \in \mathbb{Z}^n : A\mathbf{x} = \mathbf{b}\}$ .

## 1.2 Semigroups: Sparse Solutions in Integer Programming

Consider next the standard form of the feasibility constraints of integer linear programming

$$A\mathbf{x} = \mathbf{b}, \mathbf{x} \in \mathbb{Z}_{\geq 0}^n. \quad (3)$$

For a given matrix  $A$ , the set of all  $\mathbf{b}$  such that (3) is feasible, is the *semigroup*  $\mathcal{Sg}(A) = \{A\mathbf{x} : \mathbf{x} \in \mathbb{Z}_{\geq 0}^n\}$  generated by the columns of  $A$ .

If (3) has a solution, i.e.,  $\mathbf{b} \in \mathcal{Sg}(A)$ , *how sparse can such a solution be?* In other words, we are interested in the  $\ell_0$ -norm minimization problem

$$\min \{ \|\mathbf{x}\|_0 : A\mathbf{x} = \mathbf{b}, \mathbf{x} \in \mathbb{Z}_{\geq 0}^n \}. \quad (4)$$

It is clear that Problem (4) is NP-hard, because deciding the feasibility of (3) [23, §18.2] or even solving the relaxation of (4) with the condition  $\mathbf{x} \in \mathbb{Z}_{\geq 0}^n$  replaced by  $\mathbf{x} \in \mathbb{R}^n$  [19] is NP-hard.

Taking the NP-hardness of Problem (4) into account, our aim is to *estimate* the optimal value of (4) under the assumption that this problem is feasible. In [2, Theorem 1.1 (i)] (see also [1, Theorem 1]), it was shown that for any  $\mathbf{b} \in \mathcal{Sg}(A)$ , there exists a  $\mathbf{x} \in \mathbb{Z}^n$ , such that  $A\mathbf{x} = \mathbf{b}$  and

$$\|\mathbf{x}\|_0 \leq m + \left\lfloor \log_2 \left( \frac{\sqrt{\det(AA^\top)}}{\gcd(A)} \right) \right\rfloor. \quad (5)$$

In [1, Theorem 2], it was shown that (5) cannot be improved significantly, but nevertheless we show here how to improve it in some special cases. As a consequence of Theorem 1 we obtain the following.

**Corollary 3.** *Let  $A \in \mathbb{Z}^{m \times n}$  be a matrix whose columns positively span  $\mathbb{R}^m$  and let  $\mathbf{b} \in \mathbb{Z}^m$ . Then  $\mathcal{L}(A) = \mathcal{Sg}(A)$ . Furthermore, if  $\mathbf{b} \in \mathcal{L}(A)$ , and  $\tau \in \binom{[n]}{m}$  is a set, for which the matrix  $A_\tau$  is non-singular, then there is a solution  $\mathbf{x}$  of the integer-programming feasibility problem  $A\mathbf{x} = \mathbf{b}, \mathbf{x} \in \mathbb{Z}_{\geq 0}^m$  that satisfies the sparsity bound*

$$\|\mathbf{x}\|_0 \leq 2m + \Omega_m \left( \frac{|\det(A_\tau)|}{\gcd(A)} \right). \quad (6)$$

Under the above assumptions, for given  $A, \mathbf{b}$  and  $\tau$ , such a sparse solution  $\mathbf{x}$  can be computed in polynomial time.

Note that for a fixed  $m$ , (6) is usually much tighter than (5), because the function  $\Omega_m(z)$  is bounded from above by the logarithmic function  $\log_2(z)$  and is much smaller than  $\log_2(z)$  on the average. Furthermore,  $|\det(A_\tau)| \leq \sqrt{\det(AA^\top)}$  in view of the Cauchy-Binet formula.

We take a closer look at the case  $m = 1$  of a single equation and tighten the given bounds in this case. That is, we consider the *knapsack feasibility problem*

$$\mathbf{a}^\top \mathbf{x} = b, \quad \mathbf{x} \in \mathbb{Z}_{\geq 0}^n, \quad (7)$$

where  $\mathbf{a} \in \mathbb{Z}^n$  and  $b \in \mathbb{Z}$ . Without loss of generality we can assume that all components of the vector  $\mathbf{a}$  are not equal to zero. It follows from (5) that a feasible problem (7) has a solution  $\mathbf{x}$  with

$$\|\mathbf{x}\|_0 \leq 1 + \left\lfloor \log_2 \left( \frac{\|\mathbf{a}\|_2}{\gcd(\mathbf{a})} \right) \right\rfloor. \quad (8)$$

If all components of  $\mathbf{a}$  have the same sign, without loss of generality we can assume  $\mathbf{a} \in \mathbb{Z}_{>0}^n$ . In this setting, Theorem 1.2 in [2] strengthens the bound (8) by replacing the  $\ell_2$ -norm of the vector  $\mathbf{a}$  with the  $\ell_\infty$ -norm. It was conjectured in [2, page 247] that a bound  $\|\mathbf{x}\|_0 \leq c + \lfloor \log_2 (\|\mathbf{a}\|_\infty / \gcd(\mathbf{a})) \rfloor$  with an absolute constant  $c$  holds for an arbitrary  $\mathbf{a} \in \mathbb{Z}^n$ . We obtain the following result, which covers the case that has not been settled so far and yields a confirmation of this conjecture.

**Corollary 4.** *Let  $\mathbf{a} = (a_1, \dots, a_n)^\top \in (\mathbb{Z} \setminus \{0\})^n$  be a vector that contains both positive and negative components. If the knapsack feasibility problem  $\mathbf{a}^\top \mathbf{x} = b$ ,  $\mathbf{x} \in \mathbb{Z}_{\geq 0}^n$  has a solution, then there is a solution  $\mathbf{x}$  satisfying the sparsity bound*

$$\|\mathbf{x}\|_0 \leq 2 + \min \left\{ \omega \left( \frac{|a_i|}{\gcd(\mathbf{a})} \right) : i \in [n] \right\}.$$

Under the above assumptions, for given  $\mathbf{a}$  and  $b$ , such a sparse solution  $\mathbf{x}$  can be computed in polynomial time.

Our next contribution is that, given additional structure on  $A$ , we can improve on [2, Theorem 1.1 (i)], which in turn also gives an improvement on [2, Theorem 1.2]. For  $\mathbf{a}_1, \dots, \mathbf{a}_n \in \mathbb{R}^m$ , we denote by  $\text{cone}(\mathbf{a}_1, \dots, \mathbf{a}_n)$  the convex conic hull of the set  $\{\mathbf{a}_1, \dots, \mathbf{a}_n\}$ . Now assume the matrix  $A = (\mathbf{a}_1, \dots, \mathbf{a}_n) \in \mathbb{Z}^{m \times n}$  with columns  $\mathbf{a}_i$  satisfies the following conditions:

$$\mathbf{a}_1, \dots, \mathbf{a}_n \in \mathbb{Z}^m \setminus \{\mathbf{0}\}, \quad (9)$$

$$\text{cone}(\mathbf{a}_1, \dots, \mathbf{a}_n) \text{ is an } m\text{-dimensional pointed cone}, \quad (10)$$

$$\text{cone}(\mathbf{a}_1) \text{ is an extreme ray of } \text{cone}(\mathbf{a}_1, \dots, \mathbf{a}_n). \quad (11)$$

Note that the previously best sparsity bound for the general case of the integer-programming feasibility problem is (5). Using the Cauchy-Binet formula, (5) can be written as

$$\|\mathbf{x}\|_0 \leq m + \log_2 \frac{\sqrt{\sum_{I \in \binom{[n]}{m}} \det(A_I)^2}}{\gcd(A)}.$$

The following theorem improves this bound in the “*pointed cone case*” by removing a fraction of  $m/n$  of terms in the sum under the square root.

**Theorem 5.** *Let  $A = (\mathbf{a}_1, \dots, \mathbf{a}_n) \in \mathbb{Z}^{m \times n}$  satisfy (9)–(11) and, for  $\mathbf{b} \in \mathbb{Z}^m$ , consider the integer-programming feasibility problem*

$$A\mathbf{x} = \mathbf{b}, \quad \mathbf{x} \in \mathbb{Z}_{\geq 0}^n. \quad (12)$$

*If (12) is feasible, then there is a feasible solution  $\mathbf{x}$  satisfying the sparsity bound*

$$\|\mathbf{x}\|_0 \leq m + \left\lceil \log_2 \frac{q(A)}{\gcd(A)} \right\rceil,$$

where

$$q(A) := \sqrt{\sum_{I \in \binom{[n]}{m} : 1 \in I} \det(A_I)^2}.$$

We omit the proof of this result due to the page limit for the IPCO proceedings. Instead we focus on the particularly interesting case  $m = 1$ . In this case, assumption (10) is equivalent to  $\mathbf{a} \in \mathbb{Z}_{>0}^n \cup \mathbb{Z}_{<0}^n$ . Without loss of generality, one can assume  $\mathbf{a} \in \mathbb{Z}_{>0}^n$ .

**Theorem 6.** *Let  $\mathbf{a} = (a_1, \dots, a_n)^\top \in \mathbb{Z}_{>0}^n$  and  $b \in \mathbb{Z}_{\geq 0}$ . If the knapsack feasibility problem  $\mathbf{a}^\top \mathbf{x} = b$ ,  $\mathbf{x} \in \mathbb{Z}_{\geq 0}^n$  has a solution, there is a solution  $\mathbf{x}$  satisfying the sparsity bound*

$$\|\mathbf{x}\|_0 \leq 1 + \left\lceil \log_2 \left( \frac{\min\{a_1, \dots, a_n\}}{\gcd(\mathbf{a})} \right) \right\rceil.$$

When dealing with bounds for sparsity it would be interesting to understand *the worst case scenario among all members of the semigroup*, which is described by the function

$$\text{ICR}(A) = \max_{\mathbf{b} \in \mathcal{S}_g(A)} \min\{\|\mathbf{x}\|_0 : A\mathbf{x} = \mathbf{b}, \mathbf{x} \in \mathbb{Z}_{\geq 0}^n\}. \quad (13)$$

We call  $\text{ICR}(A)$  the *integer Carathéodory rank* in resemblance to the classical problem of finding the integer Carathéodory number for Hilbert bases [24]. Above results for the problem  $A\mathbf{x} = \mathbf{b}$ ,  $\mathbf{x} \in \mathbb{Z}_{\geq 0}^n$  can be phrased as upper bounds on  $\text{ICR}(A)$ . We are interested in the complexity of computing  $\text{ICR}(A)$ . The first question is: *can the integer Carathéodory rank of a matrix  $A$  be computed at all?* After all, remember that the semigroup has infinitely many elements

and, despite the fact that  $\text{ICR}(A)$  is a finite number, a direct usage of (13) would result into the determination of the sparsest representation  $A\mathbf{x} = \mathbf{b}$  for all of the infinitely many elements  $\mathbf{b}$  of  $\mathcal{Sg}(A)$ . It turns out that  $\text{ICR}(A)$  is computable, as the inequality  $\text{ICR}(A) \leq k$  can be expressed as the formula  $\forall \mathbf{x} \in \mathbb{Z}_{\geq 0}^n \exists \mathbf{y} \in \mathbb{Z}_{\geq 0}^n : (A\mathbf{x} = A\mathbf{y}) \wedge (\|\mathbf{y}\|_0 \leq k)$  in *Presburger arithmetic* [14]. Beyond this fact, the complexity status of computing  $\text{ICR}(A)$  is largely open, even when  $A$  is just one row:

**Problem 7.** *Given the input  $\mathbf{a} = (a_1, \dots, a_n)^\top \in \mathbb{Z}^n$ , is it NP-hard to compute  $\text{ICR}(\mathbf{a}^\top)$ ?*

The *Frobenius number*  $\max \mathbb{Z}_{\geq 0} \setminus \mathcal{Sg}(\mathbf{a}^\top)$ , defined under the assumptions  $\mathbf{a} \in \mathbb{Z}_{>0}^n$  and  $\gcd(\mathbf{a}) = 1$ , is yet another value associated to  $\mathcal{Sg}(\mathbf{a}^\top)$ . The Frobenius number can be computed in polynomial time when  $n$  is fixed [5, 16] but is NP-hard to compute when  $n$  is not fixed [21]. It seems that there might be a connection between computing the Frobenius number and  $\text{ICR}(\mathbf{a}^\top)$ .

## 2 Proofs of Theorem 1 and its consequences

The proof of Theorem 1 relies on the theory of finite Abelian groups. We write Abelian groups additively. An Abelian group  $G$  is said to be a *direct sum* of its finitely many subgroups  $G_1, \dots, G_m$ , which is written as  $G = \bigoplus_{i=1}^m G_i$ , if every element  $x \in G$  has a unique representation as  $x = x_1 + \dots + x_m$  with  $x_i \in G_i$  for each  $i \in [m]$ . A *primary cyclic group* is a non-zero finite cyclic group whose order is a power of a prime number. We use  $G/H$  to denote the quotient of  $G$  modulo its subgroup  $H$ .

The fundamental theorem of finite Abelian groups states that every finite Abelian group  $G$  has a *primary decomposition*, which is essentially unique. This means,  $G$  is decomposable into a direct sum of its primary cyclic groups and that this decomposition is unique up to automorphisms of  $G$ . We denote by  $\kappa(G)$  the number of direct summands in the primary decomposition of  $G$ .

For a subset  $S$  of a finite Abelian group  $G$ , we denote by  $\langle S \rangle$  the subgroup of  $G$  generated by  $S$ . We call a subset  $S$  of  $G$  *non-redundant* if the subgroups  $\langle T \rangle$  generated by proper subsets  $T$  of  $S$  are properly contained in  $\langle S \rangle$ . In other words,  $S$  is non-redundant if  $\langle S \setminus \{x\} \rangle$  is a proper subgroup of  $\langle S \rangle$  for every  $x \in S$ . The following result can be found in [13, Lemma A.6].

**Theorem 8.** *Let  $G$  be a finite Abelian group. Then the maximum cardinality of a non-redundant subset  $S$  of  $G$  is equal to  $\kappa(G)$ .*

We will also need the following lemmas, proved in the Appendix.

**Lemma 1.** *Let  $G$  be a finite Abelian group representable as a direct sum  $G = \bigoplus_{j=1}^m G_j$  of  $m \in \mathbb{Z}_{>0}$  cyclic groups. Then  $\kappa(G) \leq \Omega_m(|G|)$ .*

**Lemma 2.** *Let  $\Lambda$  be a sublattice of  $\mathbb{Z}^m$  of rank  $m \in \mathbb{Z}_{>0}^m$ . Then  $G = \mathbb{Z}^m / \Lambda$  is a finite Abelian group of order  $\det(\Lambda)$  that can be represented as a direct sum of at most  $m$  cyclic groups.*

*Proof (Theorem 1).* Let  $\mathbf{a}_1, \dots, \mathbf{a}_n$  be the columns of  $A$ . Without loss of generality, let  $\tau = [m]$ . We use the notation  $B := A_\tau$ .

*Reduction to the case  $\gcd(A) = 1$ .* For a non-singular square matrix  $M$ , the columns of  $M^{-1}A$  are representations of the columns of  $A$  in the basis of columns of  $M$ . In particular, for a matrix  $M$  whose columns form a basis of  $\mathcal{L}(A)$ , the matrix  $M^{-1}A$  is integral and the  $m \times m$  minors of  $M^{-1}A$  are the respective  $m \times m$  minors of  $A$  divided by  $\det(M) = \gcd(A)$ . Thus, replacing  $A$  by  $M^{-1}A$ , we pass from  $\mathcal{L}(A)$  to  $\mathcal{L}(M^{-1}A) = \{M^{-1}z : z \in \mathcal{L}(A)\}$ , which corresponds to a change of a coordinate system in  $\mathbb{R}^m$  and ensures that  $\gcd(A) = 1$ .

*Sparsity bound (1).* The matrix  $B$  gives rise to the lattice  $\Lambda := \mathcal{L}(B)$  of rank  $m$ , while  $\Lambda$  determines the finite Abelian group  $\mathbb{Z}^m/\Lambda$ .

Consider the canonical homomorphism  $\phi : \mathbb{Z}^m \rightarrow \mathbb{Z}^m/\Lambda$ , sending an element of  $\mathbb{Z}^m$  to its coset modulo  $\Lambda$ . Since  $\gcd(A) = 1$ , we have  $\mathcal{L}(A) = \mathbb{Z}^m$ , which implies  $\langle T \rangle = \mathbb{Z}^m/\Lambda$  for  $T := \{\phi(\mathbf{a}_{m+1}), \dots, \phi(\mathbf{a}_n)\}$ . For every non-redundant subset  $S$  of  $T$ , we have

$$\begin{aligned} |S| &\leq \kappa(\mathbb{Z}^m/\Lambda) && \text{(by Theorem 8)} \\ &\leq \Omega_m(|\det(A_\tau)|) && \text{(by Lemmas 1 and 2).} \end{aligned}$$

Fixing a set  $I \subseteq \{m+1, \dots, n\}$  that satisfies  $|I| = |S|$  and  $S = \{\phi(\mathbf{a}_i) : i \in I\}$ , we reformulate  $\langle S \rangle = \mathbb{Z}^m/\Lambda$  as  $\mathbb{Z}^m = \mathcal{L}(A_I) + \Lambda = \mathcal{L}(A_I) + \mathcal{L}(A_\tau) = \mathcal{L}(A_{I \cup \tau})$ . Thus, (1) holds for  $\gamma = I \cup \tau$ .

*Construction of  $\gamma$  in polynomial time.* The matrix  $M$  used in the reduction to the case  $\gcd(A) = 1$  can be constructed in polynomial time: one can obtain  $M$  from the Hermite Normal Form of  $A$  (with respect to the column transformations) by discarding zero columns. For the determination of  $\gamma$ , the set  $I$  that defines the non-redundant subset  $S = \{\phi(\mathbf{a}_i) : i \in I\}$  of  $\mathbb{Z}^m/\Lambda$  needs to be determined. Start with  $I = \{m+1, \dots, n\}$  and iteratively check if some of the elements  $\phi(\mathbf{a}_i) \in \mathbb{Z}^m/\Lambda$ , where  $i \in I$ , is in the group generated by the remaining elements. Suppose  $j \in I$  and we want to check if  $\phi(\mathbf{a}_j)$  is in the group generated by all  $\phi(\mathbf{a}_i)$  with  $i \in I \setminus \{j\}$ . Since  $\Lambda = \mathcal{L}(A_\tau)$ , this is equivalent to checking  $\mathbf{a}_j \in \mathcal{L}(A_{I \setminus \{j\} \cup \tau})$  and is thus reduced to solving a system of linear Diophantine equations with the left-hand side matrix  $A_{I \setminus \{j\} \cup \tau}$  and the right-hand side vector  $\mathbf{a}_j$ . Thus, carrying the above procedure for every  $j \in I$  and removing  $j$  from  $I$  whenever  $\mathbf{a}_j \in \mathcal{L}(A_{I \setminus \{j\} \cup \tau})$ , we eventually arrive at a set  $I$  that determines a non-redundant subset  $S$  of  $\mathbb{Z}^m/\Lambda$ . This is done by solving at most  $n - m$  linear Diophantine systems in total, where the matrix of each system is a sub-matrix of  $A$  and the right-hand vector of the system is a column of  $A$ .  $\square$

*Remark 1 (Optimality of the bounds).* For a given  $\Delta \in \mathbb{Z}_{\geq 2}$  let us consider matrices  $A \in \mathbb{Z}^{m \times n}$  with  $\Delta = |\det(A_\tau)|/\gcd(A)$ . We construct a matrix  $A$  that shows the optimality of the bound (1). As in the proof of Theorem 1, we assume  $\tau = [m]$  and use the notation  $B = A_\tau$ . Consider the prime factorization  $\Delta = p_1^{n_1} \cdots p_s^{n_s}$ . We will fix the matrix  $B$  to be a diagonal matrix with diagonal entries  $d_1, \dots, d_m \in \mathbb{Z}_{>0}$  so that  $\det(B) = d_1 \cdots d_m = \Delta$ .

The diagonal entries are defined by distributing the prime factors of  $\Delta$  among the diagonal entries of  $B$ . If the multiplicity  $n_i$  of the prime  $p_i$  is less than  $m$ ,

we introduce  $p_i$  as a factor of multiplicity 1 in  $n_i$  of the  $m$  diagonal entries of  $B$ . If the multiplicity  $n_i$  is at least  $m$ , we are able distribute the factors  $p_i$  among *all* of the diagonal entries of  $B$  so that each diagonal entry contains the factor  $p_i$  with multiplicity at least 1.

The group  $\mathbb{Z}^m/\Lambda = \mathbb{Z}^m/\mathcal{L}(B)$  is a direct sum of  $m$  cyclic groups  $G_1, \dots, G_m$  of orders  $d_1, \dots, d_m$ , respectively. By the Chinese Remainder Theorem, these cyclic groups can be further decomposed into the direct sum of primary cyclic groups. By our construction, the prime factor  $p_i$  of the multiplicity  $n_i < m$  generates a cyclic direct summand of order  $p_i$  in  $n_i$  of the subgroups  $G_1, \dots, G_m$ . If  $n_i \geq m$ , then each of the groups  $G_1, \dots, G_m$  has a direct summand, which is a non-trivial cyclic group whose order is a power of  $p_i$ . Summarizing, we see that the decomposition of  $\mathbb{Z}^m/\Lambda$  into primary cyclic groups contains  $n_i$  summands of order  $p_i$ , when  $n_i < m$ , and  $m$  summands, whose order is a power of  $p_i$ , when  $n_i \geq m$ . The total number of summands is thus  $\sum_{i=1}^s \min\{m, n_i\} = \Omega_m(\Delta)$ .

Now, fix  $n = m + \Omega_m(\Delta)$  and choose columns  $\mathbf{a}_{m+1}, \dots, \mathbf{a}_n$  so that  $\phi(\mathbf{a}_{m+1}), \dots, \phi(\mathbf{a}_n)$  generate all direct summands in the decomposition of  $\mathbb{Z}^m/\Lambda$  into primary cyclic groups. With this choice,  $\phi(\mathbf{a}_{m+1}), \dots, \phi(\mathbf{a}_n)$  generate  $\mathbb{Z}^m/\Lambda$ , which means that  $\mathcal{L}(A) = \mathbb{Z}^m$  and implies  $\gcd(A) = 1$ . On the other hand, any proper subset  $\{\phi(\mathbf{a}_{m+1}), \dots, \phi(\mathbf{a}_n)\}$  generates a proper subgroup of  $\mathbb{Z}^m/\Lambda$ , as some of the direct summands in the decomposition of  $\mathbb{Z}^m/\Lambda$  into primary cyclic groups will be missing. This means  $\mathcal{L}(A_{[m] \cup I}) \subsetneq \mathbb{Z}^m$  for every  $I \subsetneq \{m+1, \dots, n\}$ .

*Proof (Corollary 2).* Feasibility of (2) can be expressed as  $\mathbf{b} \in \mathcal{L}(A)$ . Choose  $\gamma$  from the assertion of Theorem 1. One has  $\mathbf{b} \in \mathcal{L}(A) = \mathcal{L}(A_\gamma)$  and so there exists a solution  $\mathbf{x}$  of (2) whose support is a subset of  $\gamma$ . This sparse solution  $\mathbf{x}$  can be computed by solving the Diophantine system with the left-hand side matrix  $A_\gamma$  and the right-hand side vector  $\mathbf{b}$ .

*Proof (Corollary 3).* Assume that the Diophantine system  $A\mathbf{x} = \mathbf{b}$ ,  $\mathbf{x} \in \mathbb{Z}^n$  has a solution. It suffices to show that, in this case, the integer-programming feasibility problem  $A\mathbf{x} = \mathbf{b}$ ,  $\mathbf{x} \in \mathbb{Z}_{\geq 0}^n$  has a solution, too, and that one can find a solution of the desired sparsity to the integer-programming feasibility problem in polynomial time.

One can determine  $\gamma$  as in Theorem 1 in polynomial time. Using  $\gamma$ , we can determine a solution  $\mathbf{x}^* = (x_1^*, \dots, x_n^*)^\top \in \mathbb{Z}^n$  of the Diophantine system  $A\mathbf{x} = \mathbf{b}$ ,  $\mathbf{x} \in \mathbb{Z}^n$  satisfying  $x_i^* = 0$  for  $i \in [n] \setminus \gamma$  in polynomial time, as described in the proof of Corollary 2.

Let  $\mathbf{a}_1, \dots, \mathbf{a}_n$  be the columns of  $A$ . Since the matrix  $A_\tau$  is non-singular, the  $m$  vectors  $\mathbf{a}_i$ , where  $i \in \tau$ , together with the vector  $\mathbf{v} = -\sum_{i \in \tau} \mathbf{a}_i$  positively span  $\mathbb{R}^n$ . Since all columns of  $A$  positive span  $\mathbb{R}^n$ , the conic version of the Carathéodory theorem implies the existence of a set  $\beta \subseteq [m]$  with  $|\beta| \leq m$ , such that  $\mathbf{v}$  is in the conic hull of  $\{\mathbf{a}_i : i \in \beta\}$ . Consequently, the set  $\{\mathbf{a}_i : i \in \beta \cup \tau\}$  and by this also the larger set  $\{\mathbf{a}_i : i \in \beta \cup \gamma\}$  positively span  $\mathbb{R}^m$ . Let  $I = \beta \cup \gamma$ . By construction,  $|I| \leq |\beta| + |\gamma| \leq m + |\gamma|$ .

Since the vectors  $\mathbf{a}_i$  with  $i \in I$  positively span  $\mathbb{R}^m$ , there exist a choice of rational coefficients  $\lambda_i > 0$  ( $i \in I$ ) with  $\sum_{i \in I} \lambda_i \mathbf{a}_i = 0$ . After rescaling we

can assume  $\lambda_i \in \mathbb{Z}_{>0}$ . Define  $\mathbf{x}' = (x'_1, \dots, x'_n)^\top \in \mathbb{Z}_{\geq 0}^n$  by setting  $x'_i = \lambda_i$  for  $i \in I$  and  $x'_i = 0$  otherwise. The vector  $\mathbf{x}'$  is a solution of  $A\mathbf{x} = \mathbf{0}$ . Choosing  $N \in \mathbb{Z}_{>0}$  large enough, we can ensure that the vector  $\mathbf{x}^* + N\mathbf{x}'$  has non-negative components. Hence,  $\mathbf{x} = \mathbf{x}^* + N\mathbf{x}'$  is a solution of the system  $A\mathbf{x} = \mathbf{b}$ ,  $\mathbf{x} \in \mathbb{Z}_{\geq 0}^n$  satisfying the desired sparsity estimate. The coefficients  $\lambda_i$  and the number  $N$  can be computed in polynomial time.

*Proof (Corollary 4).* The assertion follows by applying Corollary 3 for  $m = 1$  and all  $\tau = \{i\}$  with  $i \in [n]$ .

### 3 Proof of Theorem 6

**Lemma 3.** *Let  $a_1, \dots, a_t \in \mathbb{Z}_{>0}$ , where  $t \in \mathbb{Z}_{>0}$ . If  $t > 1 + \log_2(a_1)$ , then the system*

$$\begin{aligned} y_1 a_1 + \dots + y_t a_t &= 0, \\ y_1 \in \mathbb{Z}_{\geq 0}, \quad y_2, \dots, y_t &\in \{-1, 0, 1\}. \end{aligned}$$

*in the unknowns  $y_1, \dots, y_t$  has a solution that is not identically equal to zero.*

*Proof.* The proof is inspired by the approach in [3, §3.1] (used in a different context) that suggests to reformulate the underlying equation over integers as two strict inequalities and then use Minkowski's first theorem [4, Ch. VII, Sect. 3] from the geometry of numbers. Consider the convex set  $Y \subseteq \mathbb{R}^t$  defined by  $2t$  strict linear inequalities

$$\begin{aligned} -1 &< y_1 a_1 + \dots + y_t a_t < 1, \\ -2 &< y_i < 2 \text{ for all } i \in \{2, \dots, t\}. \end{aligned}$$

Clearly, the set  $Y$  is the interior of a hyper-parallelepiped and can also be described as  $Y = \{\mathbf{y} \in \mathbb{R}^t : \|M\mathbf{y}\|_\infty < 1\}$ , where  $M$  is the upper triangular matrix

$$M = \begin{pmatrix} a_1 & a_2 & \cdots & a_t \\ & 1/2 & & \\ & & \ddots & \\ & & & 1/2 \end{pmatrix}.$$

It is easy to see that the  $t$ -dimensional volume  $\text{vol}(Y)$  of  $Y$  is

$$\text{vol}(Y) = \text{vol}(M^{-1}[-1, 1]^t) = \frac{1}{\det(M)} 2^t = \frac{4^t}{2a_1}.$$

The assumption  $t > 1 + \log_2(a_1)$  implies that the volume of  $Y$  is strictly larger than  $2^t$ . Thus, by Minkowski's first theorem, the set  $Y$  contains a non-zero integer vector  $\mathbf{y} = (y_1, \dots, y_t)^\top \in \mathbb{Z}^t$ . Without loss of generality we can assume that  $y_1 \geq 0$  (if the latter is not true, one can replace  $\mathbf{y}$  by  $-\mathbf{y}$ ). The vector  $\mathbf{y}$  is a desired solution from the assertion of the lemma.  $\square$

*Proof (Theorem 6).* Without loss of generality we can assume that  $\gcd(\mathbf{a}) = 1$ . In fact, if  $b$  is divisible by  $\gcd(\mathbf{a})$  we can convert  $\mathbf{a}^\top \mathbf{x} = b$  to  $\bar{\mathbf{a}}^\top \mathbf{x} = \bar{b}$  with  $\bar{\mathbf{a}} = \frac{\mathbf{a}}{\gcd(\mathbf{a})}$  and  $\bar{b} = \frac{b}{\gcd(\mathbf{a})}$ , and, if  $b$  is not divisible by  $\gcd(\mathbf{a})$ , the knapsack feasibility problem  $\mathbf{a}^\top \mathbf{x} = b$ ,  $\mathbf{x} \in \mathbb{Z}_{\geq 0}^n$  has no solution.

Without loss of generality, let  $a_1 = \min\{a_1, \dots, a_n\}$ . We need to show the existence of solution of the knapsack feasibility problem satisfying  $\|\mathbf{x}\|_0 \leq 1 + \log_2(a_1)$ .

Choose a solution  $\mathbf{x} = (x_1, \dots, x_n)^\top$  of the knapsack feasibility problem with the property that the number of indices  $i \in \{2, \dots, n\}$  for which  $x_i \neq 0$  is minimized. Without loss of generality we can assume that, for some  $t \in \{2, \dots, n\}$  one has  $x_2 > 0, \dots, x_t > 0, x_{t+1} = \dots = x_n = 0$ . Lemma 3 implies  $t \leq 1 + \log_2(a_1)$ . In fact, if the latter was not true, then a solution  $\mathbf{y} \in \mathbb{R}^t$  of the system in Lemma 3 could be extended to a solution  $\mathbf{y} \in \mathbb{R}^n$  by appending zero components. It is clear that some of the components  $y_2, \dots, y_t$  are negative, because  $a_2 > 0, \dots, a_t > 0$ . It then turns out that, for an appropriate choice of  $k \in \mathbb{Z}_{\geq 0}$ , the vector  $\mathbf{x}' = (x'_1, \dots, x'_n)^\top = \mathbf{x} + k\mathbf{y}$  is a solution of the same knapsack feasibility problem satisfying  $x'_1 \geq 0, \dots, x'_t \geq 0, x'_{t+1} = \dots = x'_n = 0$  and  $x'_i = 0$  for at least one  $i \in \{2, \dots, t\}$ . Indeed, one can choose  $k$  to be the minimum among all  $a_i$  with  $i \in \{2, \dots, t\}$  and  $y_i = -1$ .

The existence of  $\mathbf{x}'$  with at most  $t - 1$  non-zero components  $x'_i$  with  $i \in \{2, \dots, n\}$  contradicts the choice of  $\mathbf{x}$  and yields the assertion.  $\square$

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## A Appendix

*Proof (Lemma 1).* Consider the prime factorization  $|G| = p_1^{n_1} \cdots p_s^{n_s}$ . Then  $|G_j| = p_1^{n_{i,j}} \cdots p_s^{n_{i,j}}$  with  $0 \leq n_{i,j} \leq n_i$  and, by the Chinese Remainder Theorem, the cyclic group  $G_j$  can be represented as  $G_j = \bigoplus_{i=1}^s G_{i,j}$ , where  $G_{i,j}$  is a cyclic group of order  $p_i^{n_{i,j}}$ . Consequently,  $G = \bigoplus_{i=1}^s \bigoplus_{j=1}^m G_{i,j}$ . This is a decomposition of  $G$  into a direct sum of primary cyclic groups and, possibly, some trivial summands  $G_{i,j}$  equal to  $\{0\}$ . We can count the non-trivial direct summands whose order is a power of  $p_i$ , for a given  $i \in [s]$ . There is at most one summand like this for each of the groups  $G_j$ . So, there are at most  $m$  non-trivial summands in the decomposition whose order is a power of  $p_i$ . On the other hand, the direct sum of all non-trivial summands whose order is a power of  $p_i$  is a group of order  $p_i^{n_{i,1} + \dots + n_{i,s}} = p_i^{n_i}$  so that the total number of such summands is not larger than  $n_i$ , as every summand contributes the factor at least  $p_i$  to the power  $p_i^{n_i}$ . This shows that the total number of non-zero summands in the decomposition of  $G$  is at most  $\sum_{i=1}^s \min\{m, n_i\} = \Omega_m(|G|)$ .  $\square$

*Proof (Lemma 2).* The proof relies on the relationship of finite Abelian groups and lattices, see [23, §4.4]. Fix a matrix  $M \in \mathbb{Z}^{m \times m}$  whose columns form a basis of  $\Lambda$ . Then  $|\det(M)| = \det(\Lambda)$ . There exist unimodular matrices  $U \in \mathbb{Z}^{m \times m}$

and  $V \in \mathbb{Z}^{m \times m}$  such that  $D := UMV$  is diagonal matrix with positive integer diagonal entries. For example, one can choose  $D$  to be the Smith Normal Form of  $M$  [23, §4.4]. Let  $d_1, \dots, d_m \in \mathbb{Z}_{>0}$  be the diagonal entries of  $D$ . Since  $U$  and  $V$  are unimodular,  $d_1 \cdots d_m = \det(D) = \det(\Lambda)$ .

We introduce the quotient group  $G' := \mathbb{Z}^m / \Lambda' = (\mathbb{Z}/d_1\mathbb{Z}) \times \cdots \times (\mathbb{Z}/d_m\mathbb{Z})$  with respect to the lattice  $\Lambda' := \mathcal{L}(D) = (d_1\mathbb{Z}) \times \cdots \times (d_m\mathbb{Z})$ . The order of  $G'$  is  $d_1 \cdots d_m = \det(D) = \det(\Lambda)$  and  $G'$  is a direct sum of at most  $m$  cyclic groups, as every  $d_i > 1$  determines a non-trivial direct summand.

To conclude the proof, it suffices to show that  $G'$  is isomorphic to  $G$ . To see this, note that  $\Lambda' = \mathcal{L}(D) = \mathcal{L}(UMV) = \mathcal{L}(UM) = \{Uz : z \in \Lambda\}$ . Thus, the map  $z \mapsto Uz$  is an automorphism of  $\mathbb{Z}^m$  and an isomorphism from  $\Lambda$  to  $\Lambda'$ . Thus,  $z \mapsto Uz$  induces an isomorphism from the group  $G = \mathbb{Z}^m / \Lambda$  to the group  $G' = \mathbb{Z}^m / \Lambda'$ .  $\square$

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