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A semigroup-theoretical view of direct-sum decompositions and associated combinatorial problems

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Let R be a ring and let $\mathcal{C}$ be a small class of right R-modules which is closed under finite direct sums, direct summands, and isomorphisms. Let $\mathcal{V}(\mathcal{C})$ denote a set of representatives of isomorphism classes in $\mathcal{C}$ and, for any module M in $\mathcal{C}$, let $[M]$ denote the unique element in $\mathcal{V}(\mathcal{C})$ isomorphic to M. Then $\mathcal{V}(\mathcal{C})$ is a reduced commutative semigroup with operation defined by $[M] + [N] = [M \oplus N]$, and this semigroup carries all information about direct-sum decompositions of modules in $\mathcal{C}$. This semigroup-theoretical point of view has been prevalent in the theory of direct-sum decompositions since it was shown

that if $\operatorname{End}_R(M)$ is semilocal for all $M \in \mathcal{C}$, then $\mathcal{V}(\mathcal{C})$ is a Krull monoid. Suppose that the monoid $\mathcal{V}(\mathcal{C})$ is Krull with a finitely generated class group (for example, when $\mathcal{C}$ is the class of finitely generated torsion-free modules and R is a one-dimensional reduced Noetherian local ring). In this case, we study the arithmetic of $\mathcal{V}(\mathcal{C})$ using new methods from zero-sum theory. Furthermore, based on module-theoretic work of Lam, Levy, Robson, and others we study the algebraic and arithmetic structure of the monoid $\mathcal{V}(\mathcal{C})$ for certain classes of modules over Prüfer rings and hereditary Noetherian prime rings.

Keywords: Krull monoids; zero-sum sequences; direct-sum decompositions; indecomposable modules; Prüfer rings; hereditary Noetherian prime rings.

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1. Introduction

The overarching goal of this manuscript is to study direct-sum decompositions of modules into indecomposable modules. Let R be a ring and let M be a right R-module. If M is Noetherian or Artinian, then a well-known and simple argument shows that M is a finite direct sum of indecomposable right R-modules. If, for example, M is either both Noetherian and Artinian or if R is a principal ideal domain and M is finitely generated, then such a direct-sum decomposition is unique; that is, the Krull–Remak–Schmidt–Azumaya property (KRSA) holds. For a simple example of non-unique direct-sum decomposition, consider a commutative domain R with distinct non-principal maximal ideals $\mathfrak{m}_1$ and $\mathfrak{m}_2$. Then the epimorphism $\mathfrak{m}_1 \oplus \mathfrak{m}_2 \to R$, defined by $(m_1, m_2) \mapsto m_1 + m_2$, gives rise to a split short exact sequence $0 \to \mathfrak{m}_1 \cap \mathfrak{m}_2 \to \mathfrak{m}_1 \oplus \mathfrak{m}_2 \to R \to 0$. Hence $\mathfrak{m}_1 \oplus \mathfrak{m}_2 \cong R \oplus (\mathfrak{m}_1 \cap \mathfrak{m}_2)$. Since $\mathfrak{m}_1$ and $\mathfrak{m}_2$ are not principal and since R is a domain, each $\mathfrak{m}_i$ is indecomposable as an R-module and is not isomorphic to R. Since the pioneering work of Krull, Remak, Schmidt, and Azumaya, direct-sum decompositions have been a classic topic in module theory. We refer the reader to [10] for an overview of the celebrated Krull–Remak–Schmidt–Azumaya Theorem and related topics in direct-sum theory.

The work of Facchini, Herbera, and Wiegand [14, 15, 17] introduced a new semigroup-theoretical approach to the study of direct-sum decompositions of modules when KRSA fails to hold. Let $\mathcal{C}$ be a small class of right R-modules which is closed under finite direct sums, direct summands, and isomorphisms. Let $\mathcal{V}(\mathcal{C})$ denote a set of representatives of isomorphism classes in $\mathcal{C}$ and, for any module M in $\mathcal{C}$, let $[M]$ denote the unique element in $\mathcal{V}(\mathcal{C})$ isomorphic to M. Then $\mathcal{V}(\mathcal{C})$ is a reduced commutative semigroup with operation defined by $[M] + [N] = [M \oplus N]$, and this semigroup carries all information about direct-sum decompositions of modules in $\mathcal{C}$. In particular, $[M]$ is an irreducible element of the semigroup $\mathcal{V}(\mathcal{C})$ if and only if M is an indecomposable module, and direct-sum decompositions of modules in $\mathcal{C}$ are unique (equivalently, KRSA holds) if and only if $\mathcal{V}(\mathcal{C})$ is a free abelian monoid. Suppose $\mathcal{C}$ is a class of R-modules as just defined and that KRSA fails.

Then arithmetical questions of the following type naturally arise.

Q1: If $M_1 \oplus M_2 = N_1 \oplus \cdots \oplus N_l$, where $M_1, M_2, N_1, \ldots, N_l$ are indecomposable modules, does there exist an upper bound for l depending only on $\mathcal{C}$?

Q2: Suppose an indecomposable module M is isomorphic to a direct summand of $N_1 \oplus \cdots \oplus N_l$ for indecomposable modules $N_1, \ldots, N_l$. Is there an upper bound (depending only on $\mathcal{C}$) for the number $|I|$ such that M is already isomorphic to a direct summand of $\bigoplus_{i \in I} N_i$?

We propose the following overall strategy to tackle these and other arithmetical questions regarding non-unique direct-sum decompositions of modules.

A. Use module-theoretic results in order to describe the algebraic structure of the semigroup $\mathcal{V}(\mathcal{C})$.

B. Use factorization theory to study the arithmetic structure of the semigroup $\mathcal{V}(\mathcal{C})$.

This strategy is relatively new, but has been used in several recent papers for certain classes of modules (see, for example [2]). In this paper, we pursue this strategy for three classes of finitely generated modules: torsion-free modules over one-dimensional reduced commutative Noetherian local rings (Sec. 4), modules over Prüfer rings (Sec. 5), and right-modules over hereditary Noetherian prime (HNP) rings (Sec. 6).

First, suppose that $\mathcal{C}$ is a class of R-modules such that the endomorphism ring $\mathrm{End}_R(M)$ is semilocal for each R-module M. In this case, Facchini proved [9, Theorem 3.4] that $\mathcal{V}(\mathcal{C})$ is a reduced Krull monoid. Earlier results in this direction can be found in [14, 15, 53]. A reduced Krull monoid H is uniquely determined by its characteristic $(G, (m_g)_{g \in G})$ where G is the class group of H and, for each $g \in G$, m_g is the cardinality of the set of prime divisors lying in the class $g \in G$ (see Sec. 2). Many arithmetical problems depend only on the set $G_p = \{g \in G : m_g > 0\}$ of classes containing prime divisors and, for simplicity, we restrict our discussion to this case. Therefore, in order to determine the structure of $\mathcal{V}(\mathcal{C})$, it is required to determine the characteristic of $\mathcal{V}(\mathcal{C})$, or at least the tuple (G, G_P). In general, this is an herculean task. Indeed, even for specific classes $\mathcal{C}$ of modules where it is known that $\mathrm{End}_R(M)$ is semilocal for each M in $\mathcal{C}$, except for in very special situations, we have limited information about (G, G_P). For an overview, what is known for certain classes of finitely generated modules over certain one- and two-dimensional Noetherian local rings, see [2]. Nevertheless, suppose we are in a situation where we are able to determine the tuple (G, G_P). Then, by a well-known transfer homomorphism (see Proposition 2.3), arithmetical problems in $\mathcal{V}(\mathcal{C})$ can be studied in the (combinatorial) Krull monoid $\mathcal{B}(G_P)$, the monoid of zero-sum sequences over G_P. Therefore, in this setting, the study of uniqueness and non-uniqueness of direct-sum decompositions can be reduced to zero-sum theory, a flourishing subfield of combinatorial and additive number theory. Except for occasional work (see [1] for

early contributions), the focus of (arithmetical) zero-sum theory has been restricted to the case where G is a finite abelian group, whereas class groups G stemming from monoids of modules $\mathcal{V}(\mathcal{C})$ are often infinite. In Sec. 3, we study zero-sum theory over finitely generated free abelian groups using new methods from matroid theory (goal **B**). The focus will be on the study of the Davenport constant which can often be used to provide bounds on important factorization-theoretic invariants. Specifically, upper and lower bounds on the Davenport constant are given in Theorem 3.13. In Sec. 4, these results will be applied to Krull monoids. In particular, module-theoretic results will be used in order to describe $\mathcal{V}(\mathcal{C})$ (goal **A**) in such a way that we can apply the results from Sec. 3 (see Corollaries 4.6 and 4.7).

Apart from cases where $\mathcal{V}(\mathcal{C})$ is Krull and some trivial cases (say, where KRSA holds whence $\mathcal{V}(\mathcal{C})$ is free abelian), the structure of $\mathcal{V}(\mathcal{C})$ has not been studied from a semigroup-theoretical point of view. Our goal in this paper is to take this approach for certain classes of finitely generated modules. The most simple case is the classical theorem of Steinitz which determines the structure of direct-sum decompositions of finitely generated modules over Dedekind domains. This result has found generalizations into many directions. We consider two generalizations in the setting of certain classes of modules over Prüfer rings and over HNP rings.

In Sec. 5, we study the class of finitely generated projective modules over a class of Prüfer rings. Based on work of Feng and Lam [18] we show that $\mathcal{V}(\mathcal{C}_{\mathrm{proj}})$ is a finitely primary monoid (and goal **A** is achieved). Since the arithmetic of finitely primary monoids has been well-studied, there are well-known answers to questions about the arithmetic of $\mathcal{V}(\mathcal{C})$ (goal **B**). These results are summarized in Theorems 5.1 and 5.3.

In Sec. 6, we study the class of finitely generated projective modules over HNP rings (which generalize noncommutative Dedekind prime rings). Deep module theoretic work by Levy and Robson [41] allows an algebraic characterization of the associated monoids of stable isomorphism classes of modules (goal **A**). We first introduce monoids of this type in an abstract setting and study their arithmetic (Propositions 6.1 and 6.4). We then apply these results to monoids of modules (goal **B**) over HNP rings in Theorem 6.5.

In both the setting of Secs. 5 and 6, the respective monoids of modules are seen to be half-factorial (all direct-sum decompositions of a given module have the same length). However, these modules behave very differently with respect to finer arithmetical invariants including the ω-invariants and the tame degrees (see Theorems 5.1, 5.3, and 6.5).

Section 2 is preparatory in nature. There we gather together arithmetical concepts from factorization theory as well as required material on (generalized) Krull monoids and monoids of modules. An even more detailed description of these concepts and their relevance to module theory can be found in [2]. We also refer the reader to the monograph [28] for more information on factorization theory, and to [3], for a friendly introduction to the interplay of factorization theory and module theory.

2. Preliminaries

We denote by $\mathbb{N}$ the set of positive integers and set $\mathbb{N}_0 = \mathbb{N} \cup \{0\}$. For real numbers $a, b \in \mathbb{R}$ we set $[a, b] = \{x \in \mathbb{Z} : a \leq x \leq b\}$. For a subset $L \subset \mathbb{Z}$, we denote by $\Delta(L) \subset \mathbb{N}$ the set of distances of L. This is the set $\{l - k : k < l \in L \text{ and } L \cap [k, l] = \{k, l\}\}$. Let G be an additive abelian group and let $A, B \subset G$ be subsets. Then $A + B = \{a + b : a \in A, b \in B\}$ denotes the sumset of A and B, $-A = \{-a : a \in A\}$ is the negative of A, $g + A = \{g\} + A$ for $g \in G$, and $\langle A \rangle \subset G$ denotes the subgroup of G generated by A.

A family $(e_i)_{i \in I}$ of elements of G is said to be *independent* if $e_i \neq 0$ for all $i \in I$ and, for every family $(m_i)_{i \in I} \in \mathbb{Z}^{(I)}$,

$$\sum_{i \in I} m_i e_i = 0 \quad \text{implies } m_i e_i = 0 \text{ for all } i \in I.$$

An independent family $(e_i)_{i \in I}$ is called a *basis* for G if $G = \bigoplus_{i \in I} \langle e_i \rangle$. If G is torsion-free, then $\mathsf{r}(G) = \dim_{\mathbb{Q}}(G \otimes_{\mathbb{Z}} \mathbb{Q})$ denotes the *rank* of G and

$$G_{|I|}^{+} = \left\{ \sum_{i \in I} \epsilon_i e_i : \epsilon_i \in \{0, 1\} \text{ for all } i \in I, \text{ but not all equal to zero} \right\}$$

$$\subset \bigoplus_{i \in I} \langle e_i \rangle \cong \mathbb{Z}^{(I)}$$

denotes the set of nonzero vertices of the hypercube in $\bigoplus_{i \in I} \langle e_i \rangle$. This definition clearly depends on the chosen basis, but throughout we refer to $G_{|I|}^{+}$ only after a basis has been fixed.

By a *monoid* we always mean a commutative semigroup with identity which satisfies the cancellation law. Thus, if R is a commutative ring and $R^{\bullet}$ its set of regular elements, then $R^{\bullet}$ is a multiplicative monoid. Let H be a (multiplicatively written) monoid. We denote by $\mathsf{q}(H)$ a quotient group of H, by $H^{\times}$ the group of invertible elements of H, and by $\widehat{H}$ the complete integral closure of H with $H \subset \widehat{H} \subset \mathsf{q}(H)$. An element $u \in H$ is called an *atom* if $u \notin H^{\times}$ and $u = ab$ with $a, b \in H$ implies that either $a \in H^{\times}$ or $b \in H^{\times}$. The set of all atoms of H is denoted by $\mathcal{A}(H)$. We say that a monoid H with identity 1 is reduced if $H^{\times} = \{1\}$, and we denote by $H_{\mathrm{red}} = \{aH^{\times} : a \in H\}$ the associated reduced monoid of H.

Free abelian monoids and groups. A monoid F is *free abelian with basis* $P \subset F$ if every $a \in F$ has a unique representation of the form

$$a = \prod_{p \in P} p^{\mathsf{v}_p(a)} \quad \text{with } \mathsf{v}_p(a) \in \mathbb{N}_0 \text{ and } \mathsf{v}_p(a) = 0 \text{ for almost all } p \in P.$$

In this case, we set $F = \mathcal{F}(P)$. The isomorphism $a \mapsto (\mathsf{v}_p(a))_{p \in P}$ from the multiplicative monoid $\mathcal{F}(P)$ to the additive monoid $(\mathbb{N}_0^{(P)}, +)$ induces an isomorphism from the quotient group $\mathsf{q}(F)$ to $\mathbb{Z}^{(P)}$. We denote by $\mathcal{F}_{\mathrm{rat}}(P)$ the multiplicative monoid isomorphic to $(\mathbb{Q}_{\geq 0}^{(P)}, +)$ and we tacitly assume that $\mathcal{F}(P) \subset \mathcal{F}_{\mathrm{rat}}(P)$. The quotient group of $\mathcal{F}_{\mathrm{rat}}(P)$ is isomorphic to $(\mathbb{Q}^{(P)}, +)$ and an element $a \in \mathsf{q}(\mathcal{F}_{\mathrm{rat}}(P))$

will be written in the form $a = \prod_{p \in P} p^{\mathsf{v}_p(a)}$ with $\mathsf{v}_p(a) \in \mathbb{Q}$ and $\mathsf{v}_p(a) = 0$ for almost all $p \in P$. We call

$$|a| = \sum_{p \in P} |\mathsf{v}_p(a)| \in \mathbb{Q}_{\geq 0} \text{ the } \textit{length} \text{ of } a \quad \text{and}$$

$$\mathrm{supp}(a) = \{p \in P : \mathsf{v}_p(a) \neq 0\} \subset P \text{ the } \textit{support} \text{ of } a.$$

Clearly there are inclusions $\mathbb{N}_0^{(P)} \subset \mathbb{Z}^{(P)} \subset \mathbb{Q}^{(P)} \subset \mathbb{Q}^P$. Elements of $\mathbb{Q}^P$ will (usually) be written as $\boldsymbol{x}$ and, if $\boldsymbol{x} \in \mathbb{Q}^P$, then we tacitly assume that $\boldsymbol{x} = (x_p)_{p \in P}$ and $\mathsf{v}_p(\boldsymbol{x}) = x_p$ for all $p \in P$. For $p \in P$, let $\boldsymbol{e}_p \in \mathbb{N}_0^{(P)}$ denote the standard vector with $e_{p,q} = 1$ if $p = q$ and $e_{p,q} = 0$ for all $q \in P \backslash \{p\}$. Then $(\boldsymbol{e}_p)_{p \in P}$ is the standard basis of $\mathbb{Z}^{(P)}$.

Factorizations and sets of lengths. Let H be a monoid. The free abelian monoid $\mathsf{Z}(H) = \mathcal{F}(\mathcal{A}(H_{\mathrm{red}}))$ is called the *factorization monoid* of H and the unique homomorphism

$$\pi : \mathsf{Z}(H) \to H_{\mathrm{red}} \quad \text{satisfying } \pi(u) = u \text{ for each } u \in \mathcal{A}(H_{\mathrm{red}})$$

is called the *factorization homomorphism* of H. For $a \in H$,

$$\mathsf{Z}_H(a) = \mathsf{Z}(a) = \pi^{-1}(aH^\times) \subset \mathsf{Z}(H) \quad \text{is the } \textit{set of factorizations} \text{ of } a,$$

$$\mathsf{L}_H(a) = \mathsf{L}(a) = \{|z| : z \in \mathsf{Z}(a)\} \subset \mathbb{N}_0 \quad \text{is the } \textit{set of lengths} \text{ of } a, \text{ and}$$

$$\mathcal{L}(H) = \{\mathsf{L}(a) : a \in H\} \quad \text{is the } \textit{system of sets of lengths} \text{ of } H.$$

We say that H is *atomic* if $\mathsf{Z}(a) \neq \emptyset$ for each $a \in H$, that H is an FF-*monoid* if $\mathsf{Z}(a)$ is finite and nonempty for each $a \in H$, and that H is *factorial* if $|\mathsf{Z}(a)| = 1$ for each $a \in H$. For the remainder of this section we assume that H is atomic.

Among the most well-studied invariants in factorization theory are those that describe the structure of sets of lengths of elements in H. Let $k \in \mathbb{N}$. If $H \neq H^\times$, then

$$\mathcal{U}_k(H) = \bigcup_{k \in L \in \mathcal{L}(H)} L$$

is the *union of sets of lengths* containing k. If $H^\times = H$, we set $\mathcal{U}_k(H) = \{k\}$. In either case we define $\rho_k(H) = \sup \mathcal{U}_k(H)$ and $\lambda_k(H) = \min \mathcal{U}_k(H)$. Clearly, $\mathcal{U}_1(H) = \{1\}$ and, for each $k \in \mathbb{N}, k \in \mathcal{U}_k(H)$. In particular, $k + l \in \mathcal{U}_k(H) + \mathcal{U}_l(H) \subset \mathcal{U}_{k+l}(H)$ for each $k, l \in \mathbb{N}$. With $\Delta(L)$ the set of distances of a length set L,

$$\Delta(H) = \bigcup_{L \in \mathcal{L}(H)} \Delta(L)$$

denotes the *set of distances* of H. By definition, $\rho_k(H) = k$ for all $k \in \mathbb{N}$ if and only if $\mathcal{U}_k(H) = \{k\}$ for all $k \in \mathbb{N}$ if and only if $\Delta(H) = \emptyset$. In this case, H is said to be *half-factorial*. If $\Delta(H) = \{d\}$ for some $d \in \mathbb{N}$, then for all $L \in \mathcal{L}(H)$ and for all $k \in \mathbb{N}, L$ and $\mathcal{U}_k(H)$ are arithmetical progressions with difference d.

Let $M \in \mathbb{N}_0, d \in \mathbb{N}$, and $\{0, d\} \subset \mathcal{D} \subset [0, d]$. A subset $L \subset \mathbb{Z}$ is called an *almost arithmetical multiprogression* (AAMP) with *difference d, period $\mathcal{D}$*, and *bound M* if

$$L = y + (L' \cup L^* \cup L'') \subset y + \mathcal{D} + d\mathbb{Z} \quad \text{where}$$

- L^* is finite and nonempty with $\min L^* = 0$ and $L^* = (\mathcal{D} + d\mathbb{Z}) \cap [0, \max L^*]$,
- $L' \subset [-M, -1]$,
- $L'' \subset \max L^* + [1, M]$, and
- $y \in \mathbb{Z}$.

Note that every AAMP is a finite nonempty subset of $\mathbb{Z}$ and that an AAMP with period $\{0, d\}$ and bound $M = 0$ is a (usual) arithmetical progression with difference d.

Distance between factorizations and catenary degrees. Let $z, z' \in \mathsf{Z}(H)$. Then we can write

$$z = u_1 \cdot \ldots \cdot u_l v_1 \cdot \ldots \cdot v_m \quad \text{and} \quad z' = u_1 \cdot \ldots \cdot u_l w_1 \cdot \ldots \cdot w_n,$$

where $l, m, n \in \mathbb{N}_0$ and $u_1, \ldots, u_l, v_1, \ldots, v_m, w_1, \ldots, w_n \in \mathcal{A}(H_{\text{red}})$ are such that

$$\{v_1, \ldots, v_m\} \cap \{w_1, \ldots, w_n\} = \emptyset.$$

Then $\gcd(z, z') = u_1 \cdot \ldots \cdot u_l$ and we call

$$\mathsf{d}(z, z') = \max\{m, n\} = \max\{|z \gcd(z, z')^{-1}|, |z' \gcd(z, z')^{-1}|\} \in \mathbb{N}_0$$

the *distance* between z and z'. The *catenary degree* $\mathsf{c}(a)$ of an element $a \in H$ is the smallest $N \in \mathbb{N}_0 \cup \{\infty\}$ such that, for any two factorizations z and z' of a, there exists a finite sequence $z = z_0, z_1, \ldots, z_k = z'$ of factorizations of a such that $\mathsf{d}(z_{i-1}, z_i) \leq N$ for all $i \in [1, k]$. We denote by

$$\mathsf{c}(H) = \sup\{\mathsf{c}(a) : a \in H\} \in \mathbb{N}_0 \cup \{\infty\}$$

the *catenary degree* of H. By definition, $|\mathsf{Z}(a)| = 1$ (i.e. a has unique factorization in H) if and only if $\mathsf{c}(a) = 0$, and thus H is factorial if and only if $\mathsf{c}(H) = 0$. Suppose now that H is not factorial. Then it is an easy consequence of the definitions that $2 + \sup \Delta(H) \leq \mathsf{c}(H)$. In particular, if $\mathsf{c}(H) = 2$, then $\Delta(H) = \emptyset$, and if $\mathsf{c}(H) = 3$, then $\Delta(H) = \{1\}$.

The ω-invariants and tame degrees. We now recall the ω-invariants as well as local and global tame degrees. These are well-studied invariants in the factorization theory of rings and semigroups. They have also been considered in module-theoretic situations in terms of the so-called *semi-exchange property* (see [6]).

For an element a in an atomic monoid H, let $\omega(H, a)$ denote the smallest $N \in \mathbb{N}_0 \cup \{\infty\}$ having the following property:

> For any multiple b of a and any factorization $b = v_1 \cdot \ldots \cdot v_n$ of b, there exists a subset $\Omega \subset [1, n]$ with $|\Omega| \leq N$ such that
> $$a \left| \prod_{\nu \in \Omega} v_\nu \right.$$

We then define

$$\omega(H) = \sup\{\omega(H, u) : u \in \mathcal{A}(H)\} \in \mathbb{N}_0 \cup \{\infty\}.$$

It is then clear that an atom $u \in H$ is prime if and only if $\omega(H, u) = 1$ and thus H is factorial if and only if $\omega(H) \leq 1$. If H satisfies the ascending chain condition (ACC) on divisorial ideals (in particular, if H is Krull or $H = R^\bullet$ where R is a Noetherian domain), then $\omega(H, u) < \infty$ for all $u \in \mathcal{A}(H)$ [29, Theorem 4.2].

Roughly speaking, for an element u in an atomic monoid H, the tame degree $\mathsf{t}(H, u)$ is the maximum of $\omega(H, u)$ and the factorization lengths of $u^{-1} \prod_{\nu \in \Omega} v_\nu$ in H. For an atom $u \in H_{\mathrm{red}}$, the *local tame degree* $\mathsf{t}(H, u)$ is the smallest $N \in \mathbb{N}_0 \cup \{\infty\}$ with the following property:

> For any multiple $a \in H_{\mathrm{red}}$ of u and any factorization $z = v_1 \cdot \ldots \cdot v_n \in \mathsf{Z}(a)$ of a which does not contain u, there is a short subproduct which is a multiple of u, say $v_1 \cdot \ldots \cdot v_m$, and a refactorization of this subproduct which contains u, say $v_1 \cdot \ldots \cdot v_m = u u_2 \cdot \ldots \cdot u_\ell$, such that $\max\{\ell, m\} \leq N$.

In particular, the local tame degree $\mathsf{t}(H, u)$ measures the distance between any factorization of a multiple a of u and a factorization of a which contains u. We denote by

$$\mathsf{t}(H) = \sup\{\mathsf{t}(H, u) : u \in \mathcal{A}(H_{\mathrm{red}})\} \in \mathbb{N}_0 \cup \{\infty\}$$

the (*global*) *tame degree* of H, and we say that H is tame if $\mathsf{t}(H) < \infty$. If u is prime, then $\mathsf{t}(H, u) = 0$ and thus H is factorial if and only if $\mathsf{t}(H) = 0$. In order to describe the relationship between the ω-invariants and the tame degree, we observe that

$$\omega(H, a) = \sup\{k \in \mathbb{N}_0 \cup \{\infty\} : b = u_1 \cdot \ldots \cdot u_k \in aH$$
$$\text{with } k \in \mathbb{N}_0, u_1, \ldots, u_k \in \mathcal{A}(H), \text{ and } u \nmid u_i^{-1} b \text{ for all } i \in [1, k]\},$$

and we define the τ-invariant as

$$\tau(H, a) = \sup\{\min \mathsf{L}(a^{-1}b) : b = u_1 \cdot \ldots \cdot u_k \in aH$$
$$\text{with } k \in \mathbb{N}_0, u_1, \ldots, u_k \in \mathcal{A}(H), \text{ and } a \nmid u_i^{-1} b \text{ for all } i \in [1, k]\}. \quad (1)$$

For each non-prime $u \in \mathcal{A}(H), \mathsf{t}(H, uH^\times) = \max\{\omega(H, u), \tau(H, u) + 1\}$ (see [29, Theorem 3.6]). Suppose that H is half-factorial and $\mathsf{L}(a) = \{l\}$ with $l \in \mathbb{N}_0$. Then $\mathsf{L}(a^{-1}b) = \{k - l\}$ and hence $\omega(H, a) = \tau(H, a) + l$. In particular, if $u \in \mathcal{A}(H)$, then $\omega(H, u) = \tau(H, u) + 1$.

We recall that if H_{red} is finitely generated, then H is tame (see [28, Theorem 3.1.4]). Tame monoids will be studied in Proposition 2.3, Theorems 5.3, and 6.5. In Sec. 4, we will characterize when the monoid $\mathcal{V}(\mathcal{C})$ is tame for $\mathcal{C}$ a class of finitely generated modules over a commutative Noetherian local ring (Proposition 4.1). Further examples of tame monoids can be found in [30, 39]. We now gather together several arithmetical finiteness properties of tame monoids.

Proposition 2.1 (Arithmetic of tame monoids). *Let H be a tame monoid.*

(1) *If H is not factorial, then $2 + \sup \Delta(H) \leq \mathsf{c}(H) \leq \omega(H) \leq \mathsf{t}(H) \leq \omega(H)^2 < \infty$.*
(2) *There is a constant $M \in \mathbb{N}_0$ such that every set of lengths $L \in \mathcal{L}(H)$ is an* AAMP *with difference $d \in \Delta(H)$ and bound M.*
(3) *There is a constant $M \in \mathbb{N}_0$ such that, for every $k \geq 2$, the union $\mathcal{U}_k(H)$ of sets of lengths is an* AAMP *with period $\{0, \min \Delta(H)\}$ and bound M.*

Proof. Statement (1) follows easily from [28, Theorem 1.6.3; 30, Sec. 3]. Statement (2) follows from [30, Theorem 5.1] and statement (3) follows from [24, Theorems 4.2 and 3.5]. □

Transfer homomorphisms. A monoid homomorphism $\theta \colon H \to B$ is called a *transfer homomorphism* if the following properties are satisfied:

(T1) $B = \theta(H)B^\times$ and $\theta^{-1}(B^\times) = H^\times$.
(T2) If $u \in H, b, c \in B$, and $\theta(u) = bc$, then there exist $v, w \in H$ such that $u = vw, \theta(v) \simeq b$, and $\theta(w) \simeq c$.

Transfer homomorphisms are a central tool in factorization theory and allow one to lift arithmetical results from a (simpler) monoid B to the monoid H (of original interest). These homomorphisms will be used throughout this manuscript (see, in particular, Propositions 2.3, 6.1, and 6.4). Each transfer homomorphism naturally gives rise to a homomorphism $\overline{\theta}$ of factorization monoids that extends θ. Let $\theta \colon H \to B$ be a transfer homomorphism of atomic monoids. Then θ induces a homomorphism $\overline{\theta} \colon \mathsf{Z}(H) \to \mathsf{Z}(B)$ satisfying $\overline{\theta}(uH^\times) = \theta(u)B^\times$ for all $u \in \mathcal{A}(H)$.

We now recall how various factorization-theoretic invariants are preserved by transfer homomorphisms. Let $\theta \colon H \to B$ and $\overline{\theta} \colon \mathsf{Z}(H) \to \mathsf{Z}(B)$ be as above. For $a \in H$, we denote by $\mathsf{c}(a, \theta)$ the smallest $N \in \mathbb{N}_0 \cup \{\infty\}$ with the following property:

If $z, z' \in \mathsf{Z}_H(a)$ and $\overline{\theta}(z) = \overline{\theta}(z')$, then there exist $k \in \mathbb{N}_0$ and factorizations $z = z_0, \ldots, z_k = z' \in \mathsf{Z}_H(a)$ such that $\overline{\theta}(z_i) = \overline{\theta}(z)$ and $\mathsf{d}(z_{i-1}, z_i) \leq N$ for all $i \in [1, k]$; that is, z and z' can be concatenated by an N-chain in the fiber $\mathsf{Z}_H(a) \cap \overline{\theta}^{-1}(\overline{\theta}(z))$.

Now

$$\mathsf{c}(H, \theta) = \sup\{\mathsf{c}(a, \theta) \colon a \in H\} \in \mathbb{N}_0 \cup \{\infty\}$$

denotes the catenary degree in the fibers of θ. The next lemma summarizes what will be needed in the sequel (a proof can be found in [28, Sec. 3.2]).

Lemma 2.2. *Let $\theta \colon H \to B$ be a transfer homomorphism.*

(1) *H is atomic if and only if B is atomic.*
(2) *For all $a \in H$, a is an atom of H if and only if $\theta(a)$ is an atom of B.*

(3) *Suppose that H is atomic.*

    (a) *For all $a \in H, \mathsf{L}_H(a) = \mathsf{L}_B(\theta(a))$. In particular,*

        (i) *$\mathcal{L}(H) = \mathcal{L}(B)$ and $\Delta(H) = \Delta(B)$, and*

        (ii) *$\mathcal{U}_k(H) = \mathcal{U}_k(B), \rho_k(H) = \rho_k(B)$, and $\lambda_k(H) = \lambda_k(B)$ for each $k \in \mathbb{N}$.*

    (b) *For all $a \in H, \mathsf{c}_B(\theta(a)) \leq \mathsf{c}_H(a) \leq \max\{\mathsf{c}_B(\theta(a)), \mathsf{c}(a, \theta)\}$. In particular,*
$$\mathsf{c}(B) \leq \mathsf{c}(H) \leq \max\{\mathsf{c}(B), \mathsf{c}(H, \theta)\}.$$

Generalized Krull monoids. Let H and D be monoids. A monoid homomorphism $\varphi : H \to D$ is called

- a *divisor homomorphism* if $\varphi(a) \,|\, \varphi(b)$ implies that $a \,|\, b$ for all $a, b \in H$.
- a *divisor theory* (for H) if $D = \mathcal{F}(P)$ for some set $P \subset D, \varphi$ is a divisor homomorphism, and, for every $a \in \mathcal{F}(P)$, there exists a finite nonempty subset $X \subset H$ such that $a = \gcd(\varphi(X))$.
- *cofinal* if for every $\alpha \in D$ there exists $a \in H$ such that $\alpha \,|\, \varphi(a)$.

If $H \subset D$ is a submonoid, then $H \subset D$ is said to be *saturated* (or *cofinal*) if the embedding $H \hookrightarrow D$ is a divisor homomorphism (or is cofinal). The monoid H is called a

- *rational generalized Krull monoid* if there is a divisor homomorphism $\varphi : H \to \mathcal{F}_{\mathrm{rat}}(P)$ for some set P.
- *Krull monoid* if there is a divisor homomorphism $\varphi : H \to \mathcal{F}(P)$ for some set P.

We note that every Krull monoid is a rational generalized Krull monoid and that a monoid H is a (rational generalized) Krull monoid if and only if H_{red} is a (rational generalized) Krull monoid. Generalized Krull monoids and domains have been studied in [4; 32, Sec. 5; 37, Chap. 22]. Specifically, [4, Proposition 2] guarantees that the definition of rational generalized Krull monoids we have given above coincides with the usual one.

    We note that a monoid is Krull if and only if it is completely integrally closed and v-noetherian and that every Krull monoid is an FF-monoid and has a divisor theory. Let H be a Krull monoid and let $\varphi : H \to D = \mathcal{F}(P)$ be a cofinal divisor homomorphism. We call $\mathcal{C}(\varphi) = \mathsf{q}(D)/\mathsf{q}(\varphi(H))$ the *class group* of φ and use additive notation for this group. For $a \in \mathsf{q}(D)$, we denote by $[a] = [a]_\varphi = a\mathsf{q}(\varphi(H)) \in \mathsf{q}(D)/\mathsf{q}(\varphi(H))$ the class containing a. Since φ is a cofinal divisor homomorphism, $\mathcal{C}(\varphi) = \{[a] : a \in D\}$ and $\varphi(H) = \{a \in D : [a] = [1]\}$. The set
$$G_P = \{[p] = p\mathsf{q}(\varphi(H)) : p \in P\} \subset \mathcal{C}(\varphi)$$

of classes containing prime divisors plays a crucial role in arithmetic computations (see Proposition 2.3) and we have $[G_P] = \mathcal{C}(\varphi)$. If φ is a divisor theory, then φ and the class group $\mathcal{C}(\varphi)$ are unique up to isomorphism. Since $\mathcal{C}(\varphi)$ depends only on H, we denote it by $\mathcal{C}(H)$ and call it the *class group of H*. Moreover, a reduced Krull monoid H with divisor theory $H \hookrightarrow \mathcal{F}(P)$ is uniquely determined

up to isomorphism by its *characteristic* $(G, (m_g)_{g \in G})$ where G is an abelian group together with an isomorphism $\Phi \colon G \to \mathcal{C}(H)$ and where $(m_g)_{g \in G}$ is a family of cardinal numbers $m_g = |P \cap \Phi(g)|$ (see [28, Theorem 2.5.4]). We consider the characteristic of a monoid of modules $\mathcal{V}(\mathcal{C})$ that is Krull in both Secs. **4** and **6**.

A domain R is a rational generalized Krull domain if and only if $R^\bullet$ is a rational generalized Krull monoid (for recent work on these kinds of domains, see [42, 47]). A v-Marot ring (and, in particular, a domain) R is a Krull ring if and only if its multiplicative monoid of regular elements $R^\bullet$ is a Krull monoid (see [33, Theorem 3.5]), and we set $\mathcal{C}(R) = \mathcal{C}(R^\bullet)$. If R is a Dedekind domain, then $\mathcal{C}(R)$ is the Picard group of R.

We now introduce a Krull monoid of combinatorial flavor, the monoid $\mathcal{B}(G_0)$ of zero-sum sequences over a subset G_0 of an abelian group G. As previously mentioned, this monoid will play a crucial role in arithmetical investigations of general Krull monoids. Section **3** provides a detailed study of $\mathcal{B}(G_0)$ in case of subsets G_0 of finitely generated free abelian groups. Let G be an additive abelian group, $G_0 \subset G$ a subset, and $S = g_1 \cdot \ldots \cdot g_l \in \mathcal{F}(G_0)$. We call $\sigma(S) = g_1 + \cdots + g_l \in G$ the sum of S, and we define, for $k \in \mathbb{N}$,

$$\Sigma(S) = \left\{ \sum_{i \in I} g_i : \emptyset \neq I \subset [1, \ell] \right\} \subset G,$$

$$\Sigma_k(S) = \left\{ \sum_{i \in I} g_i : I \subset [1, \ell], |I| = k \right\} \subset G \quad \text{and} \quad \Sigma_{\leq k}(S) = \bigcup_{\nu \in [1, k]} \Sigma_\nu(S).$$

If $(e_1, \ldots, e_r)$ is a basis of G, then the set G_r^+ of nonzero vertices of the hypercube satisfies

$$G_r^+ = \Sigma(e_1 \cdot \ldots \cdot e_r).$$

This set will be thoroughly studied in Secs. **3** and **4**.

For a map $\varphi \colon G \to G'$ between two abelian groups G and G', we define $\varphi(S) = \varphi(g_1) \cdot \ldots \cdot \varphi(g_\ell)$. Also, for $S \in \mathcal{F}(G_0)$, we set $-S = (-g_1) \cdot \ldots \cdot (-g_\ell)$. Clearly

$$\mathcal{B}(G_0) = \{S \in \mathcal{F}(G_0) : \sigma(S) = 0\} \subset \mathcal{F}(G_0)$$

is a submonoid of $\mathcal{F}(G_0)$. Moreover, since the inclusion $\mathcal{B}(G_0) \hookrightarrow \mathcal{F}(G_0)$ is a divisor homomorphism, $\mathcal{B}(G_0)$ is a Krull monoid. For arithmetical invariants $*(\cdot)$, as defined previously, we write (as it is usual) $*(G_0)$ instead of $*(\mathcal{B}(G_0))$. In particular, $\mathcal{A}(G_0)$ denotes the set of atoms of $\mathcal{B}(G_0)$, $\Delta(G_0)$ denotes the set of distances of $\mathcal{B}(G_0)$, and so forth. Note that the atoms of $\mathcal{B}(G_0)$ are precisely the minimal zero-sum sequences over G_0 — those zero-sum sequences having no proper subsequence that is also a zero-sum sequence — and we denote by

$$\mathsf{D}(G_0) = \sup\{|U| : U \in \mathcal{A}(G_0)\} \in \mathbb{N}_0 \cup \{\infty\}$$

the *Davenport constant* of G_0, a central invariant in zero-sum theory (see [25, 31, 36, 51]). The following lemma highlights the close connection between the arithmetic

of a general Krull monoid and the arithmetic of the associated monoid of zero-sum sequences. A proof can be found in [28, Theorems 3.4.2 and 3.4.10].

Proposition 2.3. *Let H be a Krull monoid, $\varphi \colon H \to D = \mathcal{F}(P)$ a cofinal divisor homomorphism, $G = \mathcal{C}(\varphi)$ its class group, and $G_P \subset G$ the set of classes containing prime divisors. Let $\widetilde{\boldsymbol{\beta}} \colon D \to \mathcal{F}(G_P)$ denote the unique homomorphism defined by $\widetilde{\boldsymbol{\beta}}(p) = [p]$ for all $p \in P$.*

(1) *The homomorphism $\boldsymbol{\beta} = \widetilde{\boldsymbol{\beta}} \circ \varphi \colon H \to \mathcal{B}(G_P)$ is a transfer homomorphism. Moreover,*
$$\mathsf{c}(H, \boldsymbol{\beta}) \leq 2, \quad \mathsf{c}(G_P) \leq \mathsf{c}(H) \leq \max\{\mathsf{c}(G_P), 2\}, \quad \text{and} \quad \mathsf{c}(H) \leq \mathsf{D}(G_P).$$

(2) *If G_P is finite, then $\mathcal{A}(G_P)$ is finite and hence $\mathsf{D}(G_P) < \infty$.*
(3) *If $\mathsf{D}(G_P) < \infty$, then both H and $\mathcal{B}(G_P)$ are tame.*

Monoids of modules. Let R be a ring and let $\mathcal{C}$ be a small class of right R-modules. That is, $\mathcal{C}$ has a set $\mathcal{V}(\mathcal{C})$ of isomorphism class representatives, and for any $M \in \mathrm{Ob}(\mathcal{C})$, we denote by $[M]$ the unique element of $\mathcal{V}(\mathcal{C})$ isomorphic to M. In more technical terms we suppose that the full subcategory $\mathcal{C}$ of Mod-R is skeletally small. Suppose that $\mathcal{C}$ is closed under finite direct sums, direct summands, and isomorphisms. Then $\mathcal{V}(\mathcal{C})$ is a reduced commutative semigroup with operation $[M] + [N] = [M \oplus N]$, and all information about direct-sum decomposition of modules in $\mathcal{C}$ can be studied in terms of factorizations in the semigroup $\mathcal{V}(\mathcal{C})$.

Suppose that $\mathcal{V}(\mathcal{C})$ is a monoid and let $\mathcal{C}'$ be a subclass of $\mathcal{C}$ which is closed under isomorphisms. Then $\mathcal{C}'$ is closed under finite direct sums, direct summands, and isomorphisms if and only if $\mathcal{V}(\mathcal{C}') \subset \mathcal{V}(\mathcal{C})$ is a divisor-closed submonoid. For a module M in $\mathcal{C}$ we denote by $\mathrm{add}(M)$ the class of R-modules that are isomorphic to direct summands of direct sums of finitely many copies of M. Then $\mathcal{V}(\mathrm{add}(M))$ is the smallest divisor-closed submonoid generated by $[M] \in \mathcal{V}(\mathcal{C})$. The class $\mathcal{C}_{\mathrm{proj}}$ of finitely generated projective right R-modules is of special importance (see [7; 13, Sec. 2.2]) and will be considered in Secs. 5 and 6.

Proposition 2.4. *Let R be a ring and $\mathcal{C}$ a small class of right R-modules which is closed under finite direct sums, direct summands, and isomorphisms.*

(1) *If $\mathrm{End}_R(M)$ is semilocal for each M in $\mathcal{C}$, then $\mathcal{V}(\mathcal{C})$ is a Krull monoid.*
(2) *If there exists M in $\mathcal{C}$ such that $\mathrm{End}_R(M)$ is semilocal, then $\mathcal{V}(\mathrm{add}(M))$ is a finitely generated Krull monoid. Conversely, if $\mathcal{V}(\mathcal{C})$ is a finitely generated monoid, then there is some M in $\mathcal{C}$ such that $\mathcal{V}(\mathcal{C}) = \mathcal{V}(\mathrm{add}(M))$.*
(3) *If R is semilocal, then $\mathcal{V}(\mathcal{C}_{\mathrm{proj}})$ is a finitely generated Krull monoid. In addition, if R is commutative, then $\mathcal{V}(\mathcal{C}_{\mathrm{proj}})$ is Krull if and only if $\mathcal{V}(\mathcal{C}_{\mathrm{proj}})$ is free abelian.*

Proof. See [9, Theorem 3.4] for the proof of (1).

We now consider statement (2). Let N_1 and N_2 be right R-modules and recall that $\mathrm{End}(N_1 \oplus N_2)$ is semilocal if and only if $\mathrm{End}(N_1)$ and $\mathrm{End}(N_2)$ are semilocal

(see the introduction of [16]). Thus $\text{End}(N)$ is semilocal for all N in $\text{add}(M)$ and hence $\mathcal{V}(\text{add}(M))$ is a Krull monoid by (1). Moreover, [8, Corollary 4.11] implies that $\mathcal{V}(\text{add}(M))$ is finitely generated. Conversely, suppose that $\mathcal{V}(\mathcal{C})$ is a finitely generated monoid. Let $M_1, \ldots, M_t$ be indecomposable right R-modules such that $\mathcal{A}(\mathcal{V}(\mathcal{C})) = \{[M_1], \ldots, [M_t]\}$. Then $M = M_1 \oplus \cdots \oplus M_t$ has the required property.

Since $\mathcal{V}(\mathcal{C}_{\text{proj}}) = \mathcal{V}(\text{add}(R_R))$ and $\text{End}(R_R) \cong R$, the first claim in (3) follows immediately from (2). The second statement follows from [12, Theorem 4.2]. $\qquad\square$

Remark 2.5. (1) Statement (1) of Proposition 2.4 is a local condition guaranteeing that $\mathcal{V}(\mathcal{C})$ is a Krull monoid and such examples will be considered in Sec. 4. It is well-known that the condition "$\text{End}_R(M)$ is semilocal" is stronger than the condition that "$\mathcal{V}(\text{add}(M))$ is a Krull monoid". Indeed, if H is any Krull monoid, then all divisor-closed submonoids are Krull as well. However, even if all divisor-closed submonoids which are generated by a single element are Krull monoids, then the monoid H can fail to be Krull (see [48] for counter-examples).

(2) In Sec. 5, we will see that since a finitely primary monoid H is Krull if and only if H is factorial, the monoid of modules $\mathcal{V}(\mathcal{C}_{\text{proj}})$ in this setting shares the same property.

(3) In Sec. 6, we will study finitely generated monoids of modules which are not Krull.

(4) Every reduced Krull monoid is isomorphic to a monoid of finitely generated projective modules $\mathcal{V}(\mathcal{C}_{\text{proj}})$ for which $\text{End}_R(M)$ is semilocal for all $M \in \mathcal{C}_{\text{proj}}$ (see [17, Theorem 2.1]). For further realization results see, for example, see [13, Sec. 9; 28, Theorem 2.7.14; 40; 53].

(5) In [11], Facchini presents an extensive list of classes of modules having semilocal endomorphism rings, and thus a list of monoids of modules $\mathcal{V}(\mathcal{C})$ which are Krull.

3. Zero-Sum Theory in Finitely Generated Free Abelian Groups

In this section, we study the Davenport constant $\mathsf{D}(G_0)$ when G_0 is a subset of a finitely generated free abelian group G. Although the results we obtain are interesting in their own right, we have in mind an application to the study of invariants of certain Krull monoids, and in particular certain monoids of modules $\mathcal{V}(\mathcal{C})$ that we study in Sec. 4.

Let G be an additive abelian group and let $G_0 \subset G$ be a subset. A sequence over G_0 will mean a finite sequence of terms from G_0 which is unordered and where repetition of terms is allowed. Zero-sum theory studies such sequences, their sets of subsequence sums, and their structure under extremal conditions (see [23, 34, 36]). Much of the focus has been on sequences over finite abelian groups, but — motivated by applications in various fields — sequences over infinite abelian groups have recently found more attention (see [26, 50] for examples of recent work).

Our goal in this section is to present, and expand upon, known methods from matroid theory that can be used to find upper bounds for the Davenport constant $\mathsf{D}(G_0)$ for subsets G_0 of finitely generated free abelian groups. These methods from matroid theory were used by Sturmfels with regards to toric varieties [52]. Such toric varieties are, in turn, related to primitive partition identities which were further studied by Diaconis, Graham, and Sturmfels [5] by combining the matroid-theoretic methods with methods from the geometry of numbers. Although no mention of either zero-sums or the Davenport constant was made in their work, Freeze and Schmid [21, Theorem 6.5] along with Geroldinger and Yuan [35, Sec. 3] made attempts to use these methods in order to study zero-sum problems over finite abelian groups. We continue this program with a goal of studying zero-sum problems over infinite abelian groups.

In order to make these methods from matroid theory — written using very different notation and terminology than is typical in zero-sum theory — more generally available for the study of zero-sum problems over finitely generated free abelian groups, we develop the basic theory here from scratch. In addition, we modify the arguments so that they extend to the case when G_0 is non-symmetric, a situation not included in the original formulations from [49, 52].

Let G be a finitely generated free abelian group, Q a $\mathbb{Q}$-vector space with $G \subset Q$, and G_0 a subset of G. The elements $S \in \mathcal{F}_{\mathrm{rat}}(G_0)$ will be called *rational sequences* over G_0. If $S = \prod_{g \in G_0} g^{\mathsf{v}_g(S)} \in \mathcal{F}_{\mathrm{rat}}(G_0)$, then $\sigma(S) = \sum_{g \in G_0} \mathsf{v}_g(S)g \in Q$ is called the *sum* of S. Then

$$\mathcal{B}_{\mathrm{rat}}(G_0) = \{S \in \mathcal{F}_{\mathrm{rat}}(G_0) : \sigma(S) = 0\} \subset \mathcal{F}_{\mathrm{rat}}(G_0)$$

is a saturated submonoid of $\mathcal{F}_{\mathrm{rat}}(G_0)$, and hence $\mathcal{B}_{\mathrm{rat}}(G_0)$ is a rational generalized Krull monoid. This monoid is clearly reduced and, since for any nonidentity $B \in \mathcal{B}_{\mathrm{rat}}(G_0)$ we have $B = B^{1/2}B^{1/2}$, $\mathcal{B}_{\mathrm{rat}}(G_0)$ has no atoms. Moreover, $\mathcal{B}_{\mathrm{rat}}(G_0)$ is a generalized block monoid as introduced in [27, Example 4.10]. The elements of $\mathcal{B}_{\mathrm{rat}}(G_0)$ will be called *rational zero-sum sequences* over G_0. Obviously, for any rational (zero-sum) sequence S there exists an integer $n \in \mathbb{N}$ such that S^n is an ordinary (zero-sum) sequence. Thus $\mathcal{B}(G_0) \subset \mathcal{B}_{\mathrm{rat}}(G_0)$ and $\mathcal{F}(G_0) \subset \mathcal{F}_{\mathrm{rat}}(G_0)$ are root extensions (see [4, Sec. 5]). If $G_0 \subset G_0'$ are two subsets of G, we assume $\mathcal{F}_{\mathrm{rat}}(G_0) \subset \mathcal{F}_{\mathrm{rat}}(G_0')$ and $\mathcal{B}_{\mathrm{rat}}(G_0) \subset \mathcal{B}_{\mathrm{rat}}(G_0')$. Specifically, we make this assumption when we consider $G_0' = G_0 \cup -G_0$. Given a sequence $S \in \mathcal{F}_{\mathrm{rat}}(G_0)$ and $g \in G_0$, we tacitly use terms such as $\mathsf{v}_{-g}(S)$ and $-S$, where we interpret $-g$ and $-S$ as elements of $\mathcal{F}_{\mathrm{rat}}(G_0 \cup -G_0)$. In particular, $\mathsf{v}_{-g}(S) = 0$ if $-g \notin G_0$.

We begin with the central construction of a partition of G_0 and of an epimorphism from $\mathcal{F}_{\mathrm{rat}}(G_0)$ to some $\mathbb{Q}$-vector space. This construction will remain valid throughout Sec. 3.

Partition the nonzero elements of G_0 as $G_0 \backslash \{0\} = G_0^+ \cup G_0^-$, where the elements of G_0 have been distributed so that if $g, -g \in G_0 \backslash \{0\}$, then g and $-g$ neither both occur in G_0^+ nor both occur in G_0^-. Furthermore, choose such a partition so that

G_0^+ is maximal, that is, such that $-G_0^- \subset G_0^+$. There may, of course, be many ways to achieve this, but we choose one such partition and fix it for the remainder of Sec. 3.

Let

$$\varphi' : \mathcal{F}_{\mathrm{rat}}(G_0) \to \mathcal{F}_{\mathrm{rat}}(G_0^+ \backslash (-G_0^-)) \times \mathsf{q}(\mathcal{F}_{\mathrm{rat}}(-G_0^-)) \subset \mathsf{q}(\mathcal{F}_{\mathrm{rat}}(G_0^+))$$

be the unique epimorphism satisfying $\varphi'(g) = g$ for each $g \in G_0^+, \varphi'(-g) = g^{-1}$ for each $-g \in G_0^-$, and $\varphi'(0) = 1$ provided $0 \in G_0$. The arguments of this section are based on the geometry of the sets and thus can be phrased more naturally using vector notation. To translate, we use the canonical isomorphism between $\mathsf{q}(\mathcal{F}_{\mathrm{rat}}(G_0^+))$ and $\mathbb{Q}^{(G_0^+)}$ which maps g to e_g for each $g \in G_0^+$ and where $(\mathsf{e}_g)_{g \in G_0^+}$ denotes the standard basis of $\mathbb{Q}^{(G_0^+)}$. Composing φ' with this isomorphism, we obtain an epimorphism

$$\varphi : \mathcal{F}_{\mathrm{rat}}(G_0) \xrightarrow{\varphi'} \mathcal{F}_{\mathrm{rat}}(G_0^+ \backslash (-G_0^-)) \times \mathsf{q}(\mathcal{F}_{\mathrm{rat}}(-G_0^-))$$

$$\to \mathbb{Q}_{\geq 0}^{(G_0^+ \backslash (-G_0^-))} \oplus \mathbb{Q}^{(-G_0^-)} \subset \mathbb{Q}^{(G_0^+)}$$

satisfying

$$\varphi \left(\prod_{g \in G_0} g^{\mathsf{v}_g(S)} \right) = \sum_{g \in G_0^+} (\mathsf{v}_g(S) - \mathsf{v}_{-g}(S)) \mathsf{e}_g.$$

The sequences $S \in \ker(\varphi) \cap \mathcal{F}(G_0)$ are precisely those zero-sum sequences over G_0 having a factorization into zero-sum subsequences of length at most 2, and an arbitrary rational sequence from $\ker(\varphi)$ is nothing more than a rational power of such a sequence. The purpose of this construction is to create a setting where we can first apply methods from linear algebra and the geometry of $\mathbb{Q}$-vector spaces to study $\varphi(S)$, and then to apply the results we obtain to the original sequence $S \in \mathcal{F}_{\mathrm{rat}}(G_0)$.

For each rational sequence $S \in \mathcal{F}_{\mathrm{rat}}(G_0)$ we define the signed support of S as

$$\mathrm{supp}^+(S) = \{g \in G_0 \cup -G_0 : \mathsf{v}_g(S) - \mathsf{v}_{-g}(S) \neq 0\}.$$

Observe that for each $S \in \mathcal{F}(G_0)$, we have

$$\mathrm{supp}^+(S) = \{g \in G_0 \cup -G_0 : \varphi'(g) \in \mathrm{supp}(\varphi'(S))\}$$

and

$$\mathrm{supp}^+(S) = \mathrm{supp}^+(-S) \subset \mathrm{supp}(S) \cup -\mathrm{supp}(S) \subset G_0 \cup -G_0.$$

For $S \in \mathcal{F}_{\mathrm{rat}}(G_0)$ we define

$$R = 0^{\mathsf{v}_0(S)} \prod_{g \in G_0^+} (g(-g))^{\min\{\mathsf{v}_g(S), \mathsf{v}_{-g}(S)\}} \in \mathcal{F}_{\mathrm{rat}}(G_0), \tag{2}$$

where $\mathsf{v}_{-g}(S) = 0$ whenever $-g \notin G_0$. Then $R \mid S$ and we set $S' = R^{-1}S \in \mathcal{F}_{\mathrm{rat}}(G_0)$. It is then easily noted that $\mathrm{supp}^+(S) = \mathrm{supp}^+(S'), \varphi(S) = \varphi(S'), \varphi(R) = \mathbf{0} \in \mathbb{Q}^{(G_0^+)}$, and for each $g \in G_0, g$ and $-g$ are not both contained in $\mathrm{supp}(S')$ for any $g \in G_0$; that is, either $\mathsf{v}_g(S') = 0$ or $\mathsf{v}_{-g}(S') = 0$. Note that the latter condition is equivalent to

$$\mathrm{supp}(S') \cap -\mathrm{supp}(S') = \emptyset.$$

In particular, $0 \notin \mathrm{supp}(S')$. Finally, if $S \in \mathcal{F}(G_0)$ is an ordinary sequence, then $R, S' \in \mathcal{F}(G_0)$ are also ordinary sequences. These observations will be used repeatedly in the sequel.

We say that a (rational) zero-sum sequence S is *elementary* if $\mathrm{supp}^+(S)$ is nonempty and minimal; that is, there is no (rational) zero-sum sequence T with $\emptyset \neq \mathrm{supp}^+(T) \subsetneq \mathrm{supp}^+(S)$. Clearly, a rational zero-sum sequence S is elementary if and only if S^n is an elementary zero-sum sequence for some positive integer n. An *elementary* atom is an atom $U \in \mathcal{B}(G_0)$ which is also an elementary zero-sum sequence. We let

$$\mathsf{D}^{\mathrm{elm}}(G_0) = \sup\{|U| : U \in \mathcal{A}(G_0) \text{ is elementary}\} \in \mathbb{N}_0 \cup \{\infty\}$$

denote the supremum of the lengths of elementary atoms over G_0 (with the convention that $\sup \emptyset = 0$), and call $\mathsf{D}^{\mathrm{elm}}(G_0)$ the *elementary Davenport constant* over G_0.

In our motivating applications the group G is the class group of a Krull monoid, and thus we need each of the previously defined concepts in this general abstract setting. However, it is technically simpler but no restriction to work over the additive group $\mathbb{Z}^r$ instead of working over an abstract finitely generated free abelian group. (Note, if $\varphi : G \to \mathbb{Z}^r$ is a group isomorphism and $G_0 \subset G$, then $\mathsf{D}(G_0) = \mathsf{D}(\varphi(G_0))$.) Therefore, for the rest of this section we suppose that

$$G_0 \subset G = \mathbb{Z}^r \quad \text{where } r \geq 1.$$

Since the case $G_0 \subset \{0\}$ is trivial, we further assume that the set G_0 contains a nonzero element of G. Moreover, whenever it is convenient, we may assume that $\mathsf{r}(\langle G_0 \rangle) = r$, as otherwise we could replace $\mathbb{Z}^r$ with $\langle G_0 \rangle \cong \mathbb{Z}^{\mathsf{r}(\langle G_0 \rangle)}$. We start with a sequence of basic but important properties regarding elementary zero-sum sequences.

Lemma 3.1. *Let $r \geq 1$ and let $G_0 \subset \mathbb{Z}^r$ be a nonempty subset. Let $S, T \in \mathcal{F}_{\mathrm{rat}}(G_0)$ be rational sequences such that $\mathrm{supp}(S) \cap \mathrm{supp}(-S) = \emptyset$ and $\mathrm{supp}(T) \cap \mathrm{supp}(-T) = \emptyset$. Suppose $\mathrm{supp}^+(T) \subset \mathrm{supp}^+(S)$ and $\mathrm{supp}(T) \cap \mathrm{supp}(S) \neq \emptyset$. Let*

$$\alpha = \min\{\mathsf{v}_g(S)/\mathsf{v}_g(T) : g \in \mathrm{supp}(S) \cap \mathrm{supp}(T)\} \in \mathbb{Q}_{>0}.$$

Define $S' = (-T)^{\alpha}S \in \mathcal{F}_{\mathrm{rat}}(G_0 \cup -G_0)$, let

$$R = \prod_{g \in G_0^+} (g(-g))^{\min\{\mathsf{v}_g(S'), \mathsf{v}_{-g}(S')\}} \in \mathcal{F}_{\mathrm{rat}}(G_0 \cup -G_0),$$

and set $\tilde{S}' := S'R^{-1}$. Then $\tilde{S}' \in \mathcal{F}_{\mathrm{rat}}(G_0)$ is a rational sequence over G_0 with

$$\mathrm{supp}(\tilde{S}') \subsetneq \mathrm{supp}(S) \quad and \quad \mathrm{supp}^+(\tilde{S}') \subsetneq \mathrm{supp}^+(S). \tag{3}$$

Proof. By definition, $\tilde{S}' \in \mathcal{B}_{\mathrm{rat}}(G_0 \cup -G_0)$ with $\mathrm{supp}^+(\tilde{S}') = \mathrm{supp}^+(S')$. Thus, by the definition of S',

$$\mathrm{supp}^+(\tilde{S}') = \mathrm{supp}^+(S') \subset \mathrm{supp}^+(-T) \cup \mathrm{supp}^+(S)$$

$$= \mathrm{supp}^+(T) \cup \mathrm{supp}^+(S) = \mathrm{supp}^+(S), \tag{4}$$

with the final equality obtained by using the hypothesis that $\mathrm{supp}^+(T) \subset \mathrm{supp}^+(S)$. By hypothesis, $\mathrm{supp}^+(T) \subset \mathrm{supp}^+(S)$, and thus

$$-g \in \mathrm{supp}(S) \subset G_0 \quad \text{for every } g \in \mathrm{supp}(T) \backslash \mathrm{supp}(S). \tag{5}$$

Let $g_1 \in \mathrm{supp}(S) \cap \mathrm{supp}(T)$ with $\mathsf{v}_g(S)/\mathsf{v}_g(T)$ minimal in $\mathbb{Q}_{>0}$ where $g \in \mathrm{supp}(S) \cap \mathrm{supp}(T)$. By hypothesis, g_1 and $-g_1$ cannot both be in $\mathrm{supp}(S)$ nor can they both be in $\mathrm{supp}(T)$. Thus, since $g_1 \in \mathrm{supp}(T) \cap \mathrm{supp}(S)$, we see that $-g_1 \notin \mathrm{supp}(T)$ and $-g_1 \notin \mathrm{supp}(S)$. Consequently, by definition of α, we have

$$\mathsf{v}_{-g_1}(-T^\alpha) = \mathsf{v}_{g_1}(T^\alpha) = \mathsf{v}_{g_1}(S) \quad and \quad \mathsf{v}_{g_1}(-T^\alpha) = \alpha\mathsf{v}_{-g_1}(T) = 0 = \mathsf{v}_{-g_1}(S).$$

But $g_1 \in \mathrm{supp}^+(S)$ as $g_1 \in \mathrm{supp}(S)$ and $-g_1 \notin \mathrm{supp}(S)$. Thus $g_1 \notin \mathrm{supp}^+(S')$ and the inclusion in (4) must be strict; that is,

$$\mathrm{supp}^+(\tilde{S}') \subsetneq \mathrm{supp}^+(S). \tag{6}$$

From the definitions of S' and $\tilde{S}'$, we have

$$\mathrm{supp}(\tilde{S}') \subset \mathrm{supp}(S') \subset \mathrm{supp}(-T) \cup \mathrm{supp}(S). \tag{7}$$

For $g \in \mathrm{supp}(T) \cap \mathrm{supp}(S)$, the definition of α ensures that $\mathsf{v}_{-g}(-T^\alpha) \leq \mathsf{v}_g(S)$. Since g and $-g$ cannot both be in $\mathrm{supp}(T)$ nor both be in $\mathrm{supp}(S)$, we also have that $\mathsf{v}_g(-T^\alpha) = 0 = \mathsf{v}_{-g}(S)$. Combining this fact with the definition of $\tilde{S}'$ we see that for every $g \in \mathrm{supp}(T) \cap \mathrm{supp}(S)$, $-g \notin \mathrm{supp}(\tilde{S}')$. Considering this fact along with (5) and (7), it follows that

$$\mathrm{supp}(\tilde{S}') \subset \mathrm{supp}(S) \subset G_0. \tag{8}$$

In particular, we now know that $\tilde{S}' \in \mathcal{F}_{\mathrm{rat}}(G_0)$. It remains to show that the inclusion in (8) is strict. But for this fact we need to only recall that $g_1 \in \mathrm{supp}(S)$ and that $g_1 \notin \mathrm{supp}^+(\tilde{S}') = \mathrm{supp}(\tilde{S}') \cup -\mathrm{supp}(\tilde{S}')$. Indeed, this last equality follows immediately since $\mathrm{supp}(\tilde{S}')$ contains at most one of g and $-g$ for every $g \in G_0$, a fact that follows from the definition of $\mathrm{supp}(\tilde{S}')$ and the observations before Lemma 3.1. □

Lemma 3.2. *Let $r \geq 1$ and let $G_0 \subset \mathbb{Z}^r$ be a nonempty subset. If $S, T \in \mathcal{B}_{\mathrm{rat}}(G_0)$ are both elementary rational zero-sum sequences with $\mathrm{supp}^+(S) = \mathrm{supp}^+(T)$, then either*

$$SR_T^\alpha = T^\alpha R_S \quad or \quad S(-R_T)^\alpha = -T^\alpha R_S \quad for some positive \alpha \in \mathbb{Q}_{>0},$$

where $R_S \mid S$ and $R_T \mid T$ are the respective maximal length rational zero-sum subsequences of S or T with $\varphi(R_T) = 0$ and $\varphi(R_S) = 0$ (defined in (2)).

In particular, if $S, T \in \mathcal{B}(G_0)$ are elementary zero-sum sequences with common signed support, then there exist relatively prime $m, n \in \mathbb{N}$ such that either

$$S^n R_T^m = T^m R_S^n \quad or \quad S^n(-R_T)^m = (-T)^m R_S^n,$$

where $R_S \mid S$ and $R_T \mid T$ are the respective maximal length zero-sum subsequences of S or T having a factorization into zero-sum subsequences each of length at most 2.

Proof. Let $S, T \in \mathcal{B}_{\mathrm{rat}}(G_0)$ be elementary rational zero-sum sequences with common signed support $\mathrm{supp}^+(S) = \mathrm{supp}^+(T)$. In view of the observations before Lemma 3.1 we may, without loss of generality, assume that R_S and R_T are trivial (the general case follows easily from this special case). In particular,

$$\mathrm{supp}(S) \cap \mathrm{supp}(-S) = \emptyset \quad and \quad \mathrm{supp}(T) \cap \mathrm{supp}(-T) = \emptyset. \tag{9}$$

Moreover, since S and T are elementary zero-sum sequences, $\mathrm{supp}^+(S) = \mathrm{supp}^+(T) \neq \emptyset$.

First suppose that $\mathrm{supp}(S) \cap \mathrm{supp}(T) \neq \emptyset$. Then, since $\mathrm{supp}^+(S) = \mathrm{supp}^+(T)$, we can apply Lemma 3.1. Let $\alpha \in \mathbb{Q}_{>0}$, $S' \in \mathcal{F}_{\mathrm{rat}}(G_0 \cup -G_0)$, and $\tilde{S}' \in \mathcal{F}_{\mathrm{rat}}(G_0)$ be as in Lemma 3.1. Since S, T, and R are each rational zero-sum sequences, it follows from the definition of $\tilde{S}'$ that $\tilde{S}'$ is also a rational zero-sum sequence. Thus $\tilde{S}'$ will, in view of (3), contradict that S is an elementary rational zero-sum sequence unless $\mathrm{supp}^+(\tilde{S}') = \emptyset$. However, by definition, $\mathrm{supp}^+(\tilde{S}') = \emptyset$ is only possible if $\tilde{S}'$ is trivial, in which case $(-T^\alpha)S = S' = R$. But this implies, in view of (9) and $\mathrm{supp}^+(T) = \mathrm{supp}^+(S)$, that

$$\mathsf{v}_g(T^\alpha) = \mathsf{v}_g(S) \quad \text{for all } g \in G_0.$$

Therefore, $S = T^\alpha$ as desired.

We now assume that $\mathrm{supp}(S) \cap \mathrm{supp}(T) = \emptyset$. In this case, since $\mathrm{supp}^+(S) = \mathrm{supp}^+(T)$, it follows from (9) that $\mathrm{supp}(T) = -\mathrm{supp}(S)$. This in turn implies that $\mathrm{supp}(-T) \subset G_0$ and hence $-T \in \mathcal{B}_{\mathrm{rat}}(G_0)$. Repeating the arguments of the previous paragraph using $-T$ in place of T, we conclude that $S = (-T)^\alpha$ as desired. Then, in particular, statement follows easily from the general statement with $\alpha = \frac{m}{n}$, where $m, n \in \mathbb{N}$ are relatively prime. $\qquad\square$

Lemma 3.3. *Let $r \geq 1$ and let $G_0 \subset \mathbb{Z}^r$ be a nonempty subset. If $U, V \in \mathcal{A}(G_0)$ are elementary atoms with $\mathrm{supp}^+(U) = \mathrm{supp}^+(V)$, then either $U = V$ or $U = -V$.*

Proof. Since U and V are both elementary, their signed support must be nontrivial. Moreover, since U and V are atoms, neither U nor V has a nontrivial zero-sum subsequence of length one or two. Therefore, applying Lemma 3.2 to both U and V, we find that either $U^n = V^m$ or $U^n = (-V)^m$ for relatively prime positive integers m and n. Note that if the latter case holds, then $\mathrm{supp}(-V) = \mathrm{supp}(U) \subset G_0$.

Thus, by replacing V with $-V$ if need be (in which case the hypotheses of the theorem hold for U and $-V$ and, if the conclusion holds for U and $-V$, then it will hold for the original pair U and V as well), we may assume that $U^n = V^m$ for relatively prime positive integers m and n. Therefore, $n\mathsf{v}_g(U) = m\mathsf{v}_g(V)$ for all $g \in G_0$. Consequently, since $\gcd(m, n) = 1$, it follows that $m \mid \mathsf{v}_g(U)$ for each $g \in G_0$. But then $U = W^m$ is a product decomposition of U, where $W = \prod_{g \in G_0} g^{\mathsf{v}_g(U)/m} \in \mathcal{F}(G_0)$. Noting that $0 = \sigma(U) = m\sigma(W)$, we conclude that W is a zero-sum sequence. Since it was assumed that $U \in \mathcal{A}(G_0)$ is an atom, $m = 1$. A similar argument shows that $n = 1$. Now the relation $U^m = V^m$ gives the desired conclusion $U = V$. $\qquad\square$

Lemma 3.4. *Let $r \geq 1$ and let $G_0 \subset \mathbb{Z}^r$ be a nonempty subset. If $S \in \mathcal{B}(G_0)$ is an elementary zero-sum sequence, then*

$$S = RU^\ell$$

for some $\ell \geq 1$, some elementary atom $U \in \mathcal{A}(G_0)$ with $\mathrm{supp}^+(U) = \mathrm{supp}^+(S)$, and some zero-sum sequence $R \in \ker(\varphi)$ that has a factorization involving only zero-sum subsequences all of length at most 2.

Proof. Assume for the sake of contradiction that $S \in \mathcal{B}(G_0)$ is an elementary zero-sum sequence that fails to have the desired form and with $|S|$ minimal. By the minimality of $|S|$ and the observations made before Lemma 3.1, it follows that every nontrivial zero-sum subsequence of S has length at least 3. Therefore $\mathrm{supp}(S) \cap \mathrm{supp}(-S) = \emptyset$. Since the conclusion of the lemma holds for each atom, it must be the case that the chosen S is not an atom. Therefore,

$$S = T_1 T_2$$

for nontrivial zero-sum subsequences T_1 and T_2. Now, since every zero-sum subsequence of S has length at least 3, the nontrivial zero-sum subsequences T_1 and T_2 must each have nonempty signed support. But this contradicts the fact that S is an elementary zero-sum sequence unless we have

$$\mathrm{supp}^+(S) = \mathrm{supp}^+(T_1) = \mathrm{supp}^+(T_2). \tag{10}$$

Also, $|T_1| < |S|$ and $|T_2| < |S|$ since T_1 and T_2 are both nontrivial subsequences of S. By (10), each T_i is a zero-sum sequence with $\mathrm{supp}^+(T_i) = \mathrm{supp}^+(S)$ and is consequently an elementary zero-sum sequence. Furthermore, since $|T_1| < |S|$ and $|T_2| < |S|$, the minimality of $|S|$ ensures that both T_1 and T_2 have the form stated in the conclusion of the lemma. Thus $T_1 = U^m$ and $T_2 = V^n$ for atoms $U, V \in \mathcal{A}(G_0)$ with

$$\mathrm{supp}^+(U) = \mathrm{supp}^+(V) = \mathrm{supp}^+(T_1) = \mathrm{supp}^+(T_2) = \mathrm{supp}^+(S)$$

and positive integers m and n. As U and V are both atoms with $\mathrm{supp}^+(U) = \mathrm{supp}^+(V) = \mathrm{supp}^+(S)$, it follows that U and V are both elementary atoms. Invoking Lemma 3.3 we find that either $U = V$ or $U = -V$.

If $U = -V$, then there must exist $g \in \mathrm{supp}(U)$ with $-g \in \mathrm{supp}(V)$. In this case, since $T_1 = U^m, T_2 = V^n$, and $S = T_1 T_2$, it follows that $g(-g) \mid S$, contrary to conclusion above that S does not have any zero-sum subsequence of length 1 or 2. Therefore $U \neq -V$, forcing $U = V$. Then $T_1 = U^m, T_2 = V^n = U^n$, and $S = T_1 T_2 = U^{m+n}$ as desired. $\qquad\square$

Combining the above results, we are now able to describe the form of elementary zero-sum sequences. We do so in the following proposition whose proof follows immediately from Lemmas 3.3 and 3.4. Essentially, Proposition 3.5 states that if X is the signed support of an elementary zero-sum sequence, then (up to sign) there is a unique atom U with $\mathrm{supp}^+(U) = X$, and all other zero-sum sequences having signed support X must have the form $U^\ell R$ or $(-U)^\ell R$ where $\ell \geq 1$ and $R \in \mathcal{B}(G_0)$ is a zero-sum sequence that has a factorization involving zero-sum subsequences all of length at most 2.

Proposition 3.5. *Let $r \geq 1$ and let $G_0 \subset \mathbb{Z}^r$ be a nonempty subset. If X is the signed support of an elementary zero-sum sequence over G_0, then there exists a unique (up to sign) atom $U \in \mathcal{A}(G_0)$ with $\mathrm{supp}^+(U) = \mathrm{supp}^+(-U) = X$ such that every elementary zero-sum sequence S with signed support $\mathrm{supp}^+(S) = X$ has the form*

$$S = RU^\ell \quad or \quad S = R(-U)^\ell,$$

where $R \in \ker(\varphi)$ is a zero-sum sequence and $\ell \geq 1$.

Lemma 3.6. *Let $r \geq 1$ and let $G_0 \subset \mathbb{Z}^r$ be a nonempty subset. Suppose $S \in \mathcal{B}_{\mathrm{rat}}(G_0)$ is a rational zero-sum sequence with $\mathrm{supp}^+(S)$ nonempty. Then there exists some elementary atom $U \in \mathcal{A}(G_0)$ such that $\mathrm{supp}(U) \subset \mathrm{supp}(S)$.*

Proof. Assume for the sake of contradiction that $S \in \mathcal{B}_{\mathrm{rat}}(G_0)$ is a counterexample with $\mathrm{supp}^+(S) \neq \emptyset$ minimal. By removing rational zero-sum subsequences of the form $((-g)g)^{\min\{\mathsf{v}_g(S), \mathsf{v}_{-g}(S)\}}$ and $0^{\mathsf{v}_0(S)}$ from S as defined in (2), we may, without loss of generality, assume that $\mathrm{supp}(S) \cap \mathrm{supp}(-S) = \emptyset$. Let $U \in \mathcal{B}(G_0)$ be an elementary zero-sum sequence with

$$\mathrm{supp}^+(U) \subset \mathrm{supp}^+(S).$$

Note that such an elementary zero-sum sequence exists with $\mathrm{supp}^+(U) = X$ for any minimal nonempty subset $X \subset \mathrm{supp}^+(S)$ provided there exists a zero-sum sequence with signed support X.

In view of the observations made prior to Lemma 3.1, we may assume that U has no nontrivial zero-sum sequence of length 1 or 2. Now, if $U = U_1 \cdot \ldots \cdot U_\ell$ is a factorization as a product of atoms $U_i \in \mathcal{A}(G_0)$, we find that $\mathrm{supp}^+(U_i) = \mathrm{supp}^+(U)$ for each $i \in [1, l]$, else U is not an elementary zero-sum sequence. Therefore, replacing U by some U_i as need be, we may without loss of generality assume that $U \in \mathcal{A}(G_0)$ is an atom.

Since $\mathrm{supp}^+(U) \subset \mathrm{supp}^+(S), \mathrm{supp}(U)$ and $\mathrm{supp}(S)$ are disjoint only if $-g \in \mathrm{supp}(S) \subset G_0$ for every $g \in \mathrm{supp}(U)$. In this case, $-U \in \mathcal{A}(G_0)$ is also an atom over G_0. Thus, replacing U by $-U$ (in this one scenario), we may assume that $\mathrm{supp}(U) \cap \mathrm{supp}(S) \neq \emptyset$. Since an elementary atom cannot have a nontrivial zero-sum sequence of length 1 or 2, we may apply Lemma 3.1 with $T = U$. Now let α, S', R and $\tilde{S}'$ be as in Lemma 3.1. Since S, U and R are each zero-sum sequences, it follows from the definition of $\tilde{S}'$ that $\tilde{S}'$ is also a zero-sum sequence. Therefore,

$$\tilde{S}' \in \mathcal{B}_{\mathrm{rat}}(G_0),$$

an improvement over $\tilde{S}' \in \mathcal{F}_{\mathrm{rat}}(G_0)$ given in Lemma 3.1.

If $\mathrm{supp}^+(\tilde{S}') = \mathrm{supp}^+(S')$ is nonempty, then by the strict inclusion $\mathrm{supp}^+(\tilde{S}') \subsetneq \mathrm{supp}^+(S)$ in (3), we may apply the induction hypothesis to $\tilde{S}' \in \mathcal{B}_{\mathrm{rat}}(G_0)$ and find an elementary atom $V \in \mathcal{A}(G_0)$ such that $\mathrm{supp}(V) \subset \mathrm{supp}(\tilde{S}')$. The proof is then complete if one considers the other inclusion in (3). Therefore, we may assume the alternative, that

$$\mathrm{supp}^+(\tilde{S}') = \mathrm{supp}^+(S') = \emptyset. \tag{11}$$

Recalling that $S' = (-U)^\alpha S$, that $\mathrm{supp}^+(U) \subset \mathrm{supp}^+(S)$, and that $\mathrm{supp}(S) \cap \mathrm{supp}(-S) = \mathrm{supp}(U) \cap \mathrm{supp}(-U) = \emptyset$, we see that (11) is possible only if $\mathrm{supp}(U) = \mathrm{supp}(S)$ with $\mathsf{v}_g(U^\alpha) = \mathsf{v}_g(S)$ for every $g \in \mathrm{supp}(U) = \mathrm{supp}(S)$. $\qquad\square$

Let $G_0 \subset \mathbb{Z}^r \subset \mathbb{Q}^r$ be a finite subset. Note that the chosen partition of $G_0 \backslash \{0\}$ gives rise to a unique partition of $G_0 \cup -G_0 \backslash \{0\}$ with $(G_0 \cup -G_0)^+ = G_0^+$. To this partition we may again associate a map $\mathcal{F}_{\mathrm{rat}}(G_0 \cup -G_0) \to \mathbb{Q}^{(G_0^+)}$ which we also denote by φ. Then $\varphi : \mathcal{F}_{\mathrm{rat}}(G_0) \to \mathbb{Q}^{(G_0^+)}$ is simply the restriction of $\varphi : \mathcal{F}_{\mathrm{rat}}(G_0 \cup -G_0) \to \mathbb{Q}^{(G_0^+)}$ to $\mathcal{F}_{\mathrm{rat}}(G_0)$. By construction, we have

$$\varphi(\mathcal{B}(G_0)) \subset \mathbb{Z}^{(G_0^+)} \subset \mathbb{Q}^{(G_0^+)}.$$

It is easily checked that $\varphi(\mathcal{B}_{\mathrm{rat}}(G_0))$ is an additive monoid closed under multiplication by nonnegative rational numbers. The $\mathbb{Q}$-vector space spanned by $\varphi(\mathcal{B}_{\mathrm{rat}}(G_0))$ is then $\varphi(\mathcal{B}_{\mathrm{rat}}(G_0 \cup -G_0))$, which is also the $\mathbb{Q}$-vector space spanned by $\varphi(\mathcal{B}(G_0)) = \varphi(\mathcal{B}_{\mathrm{rat}}(G_0)) \cap \mathbb{Z}^{(G_0^+)}$.

Note that vector $(\alpha_g)_{g \in G_0^+} \in \mathbb{Q}^{(G_0^+)}$ is an element of $\varphi(\mathcal{B}_{\mathrm{rat}}(G_0 \cup -G_0))$ precisely when

$$\sum_{g \in G_0^+} \alpha_g g = 0.$$

Thus, if we let M denote the $r \times |G_0^+|$ matrix whose columns are the vectors $g \in G_0^+ \subset \mathbb{Z}^r \subset \mathbb{Q}^r$, we see that $\varphi(\mathcal{B}_{\mathrm{rat}}(G_0 \cup -G_0))$ is the kernel of the matrix M. The set $\varphi(\mathcal{B}_{\mathrm{rat}}(G_0))$ can also be described via M; It is the subset consisting of all vectors $(\alpha_g)_{g \in G_0^+} \in \ker(M)$ that satisfy the following sign restrictions:

$$\alpha_g \geq 0 \quad \text{unless } -g \in G_0.$$

(Recall that we always assume $-G_0^- \subset G_0^+$, and thus $\alpha_g > 0$ is allowed for every $g \in G_0^+$.) It is well-known that the kernel of a matrix M is the orthogonal space for the row space of the same matrix M, and that the row and column space of M have the same dimension, which in this case is equal to the dimension of the $\mathbb{Q}$-vector space spanned by the vectors from $G_0 \subset \mathbb{Z}^r \subset \mathbb{Q}^r$. This latter number is simply $\mathsf{r}(\langle G_0 \rangle)$ and thus we conclude that

$$|G_0^+| = \mathsf{r}(\langle G_0 \rangle) + \dim_{\mathbb{Q}}(\langle \varphi(\mathcal{B}(G_0)) \rangle). \tag{12}$$

The next theorem (essentially due to Rockafellar [49] in a matroid formulation) shows that elementary zero-sum sequences can be useful for decomposing a zero-sum sequence via rational product decomposition. Indeed, an arbitrary zero-sum sequence always has a product decomposition into a bounded number of rational powers of elementary atoms. It also shows that an upper bound for $\mathsf{D}(G_0)$ can be found using an upper bound for $\mathsf{D}^{\mathsf{elm}}(G_0)$. It is important to note that, even if $S \in \mathcal{B}(G_0)$ is an atom, S may still have a nontrivial product decomposition into rational powers of elementary atoms if it is not itself elementary.

Theorem 3.7. *Let $r \geq 1$, let $G_0 \subset \mathbb{Z}^r$ be a nonempty subset, and let $S \in \mathcal{B}_{\mathrm{rat}}(G_0)$ be a rational zero-sum sequence. Then there exist a nonnegative integer $\ell \geq 0$, elementary atoms $U_1, \ldots, U_\ell \in \mathcal{A}(G_0)$, positive rational numbers $\alpha_0, \alpha_1, \ldots, \alpha_\ell \in \mathbb{Q}_{>0}$, and a zero-sum sequence $R \in \ker(\varphi)$ (possibly trivial) having a factorization into zero-sum subsequences each of length 1 or 2. Moreover,*

$$S = R^{\alpha_0} U_1^{\alpha_1} \cdot \ldots \cdot U_\ell^{\alpha_\ell}, \quad \mathrm{supp}^+(U_j) \nsubseteq \mathrm{supp}^+ \left(S \left(\prod_{i=1}^j U_i^{\alpha_i} \right)^{-1} \right)$$

for all $i \in [1, \ell]$, and

$$\ell \leq \min \left\{ \frac{1}{2} |\mathrm{supp}^+(S)|, |G_0^+| - \mathsf{r}(\langle G_0 \rangle) \right\}.$$

In particular,

$$\mathsf{D}(G_0) \leq \sup\{2, \hat{\ell}\, \mathsf{D}^{\mathsf{elm}}(G_0)\} \leq \sup\{2, \min\{\eta, |G_0^+| - \mathsf{r}(\langle G_0 \rangle)\} \mathsf{D}^{\mathsf{elm}}(G_0)\}$$

$$\leq \sup\{2, |G_0 \backslash \{0\}| \mathsf{D}^{\mathsf{elm}}(G_0)\},$$

where $\eta = \sup\{|\mathrm{supp}(U)| : U \in \mathcal{A}(G_0)\}$ and $\hat{\ell}$ denotes the supremum over all ℓ as S ranges over $\mathcal{A}(G_0)$.

Proof. We first construct the rational product decomposition for $S \in \mathcal{B}_{\mathrm{rat}}(G_0)$ and then deduce the upper bound for $\mathsf{D}(G_0)$. To this end, in view of the observations made prior to Lemma 3.1, we may without loss of generality assume that S is nontrivial and that

$$\{g \in G_0^+ : g, -g \in \mathrm{supp}(S)\} = \emptyset. \tag{13}$$

It suffices to prove the theorem for such rational sequences.

Our goal is to show that there exist elementary atoms $U_1, \ldots, U_\ell \in \mathcal{A}(G_0)$ and positive rational numbers $\alpha_1, \ldots, \alpha_\ell \in \mathbb{Q}_{>0}$ such that

$$S = U_1^{\alpha_1} \cdot \ldots \cdot U_\ell^{\alpha_\ell} \quad \text{and} \quad \mathrm{supp}^+(U_j) \nsubseteq \mathrm{supp}^+ \left(S \left(\prod_{i=1}^{j} U_i^{\alpha_i} \right)^{-1} \right)$$

$$\text{for all } i \in [1, \ell]. \qquad (14)$$

Before doing so, we explain how (14) forces the desired upper bound for ℓ. The second condition from (14) says that each atom U_j contains an element $g_j \in \mathrm{supp}(S)$ (via the first condition of (14)) not contained in any of the $U_{j+1}, \ldots, U_\ell$. But then $\{g_1, \ldots, g_\ell\} \subset \mathrm{supp}(S)$ is a subset of cardinality ℓ, implying $\ell \leq |\mathrm{supp}(S)| = \frac{1}{2}|\mathrm{supp}^+(S)|$. In view of (13) and since $g_j \in \mathrm{supp}(S)$, we see that $-g_j \notin \mathrm{supp}(S)$ for each $j \in [1, \ell]$. As a result, since each atom U_j contains some element g_j not contained in any of the $U_{j+1}, \ldots, U_\ell$, no $g_j \in \mathrm{supp}^+(U_j)$ is contained in $\mathrm{supp}^+(U_{j+1} \cdot \ldots \cdot U_\ell)$. Now it is easily deduced that $\varphi(U_1), \ldots, \varphi(U_\ell) \in \varphi(\mathcal{B}(G_0))$ are linearly independent over $\mathbb{Q}$. Therefore, $\ell \leq \dim_{\mathbb{Q}}(\langle \varphi(\mathcal{B}(G_0)) \rangle) = |G_0^+| - \mathsf{r}(\langle G_0 \rangle)$, with the equality following from (12). Thus the desired bound for ℓ follows from (14), and we now devote our attention to proving (14). For this, suppose for the sake of contradiction that $S \in \mathcal{B}_{\mathrm{rat}}(G_0)$ is a counterexample to (14) with $|\mathrm{supp}(S)|$ minimal.

Since S is nontrivial, (13) forces $\mathrm{supp}^+(S) \neq \emptyset$. Thus, by Lemma 3.6, there exists an elementary atom $U_1 \in \mathcal{A}(G_0)$ with $\mathrm{supp}(U_1) \subset \mathrm{supp}(S)$. Let $\alpha_1 = \min\{\mathsf{v}_g(S)/\mathsf{v}_g(U_1) : g \in \mathrm{supp}(U_1)\}$. Since $\mathrm{supp}(U_1) \subset \mathrm{supp}(S), \alpha_1 > 0$. By definition, $\mathsf{v}_g(U_1^{\alpha_1}) \leq \mathsf{v}_g(S)$ for every $g \in \mathrm{supp}(U_1)$, with equality holding for some $g = g_1 \in \mathrm{supp}(U_1) \subset \mathrm{supp}(S)$ (attaining the minimum in the definition of α_1). Define $S' = SU_1^{-\alpha_1}$. Since $\mathsf{v}_{g_1}(U_1^{\alpha_1}) = \mathsf{v}_{g_1}(S)$ with $g_1 \in \mathrm{supp}(S)$, we conclude that $\mathrm{supp}(S')$ is a proper subset of $\mathrm{supp}(S)$. Indeed, $g_1 \notin \mathrm{supp}(S')$. Since S and U_1 are each zero-sum sequences, it follows that S' is also a zero-sum sequence. Thus, in view of the minimality of $|\mathrm{supp}(S)|$, we can apply the theorem to (the possibly trivial) rational sequence $\tilde{S}'$ to find $\tilde{S}' = U_2^{\alpha_2} \cdot \ldots \cdot U_\ell^{\alpha_\ell}$ for some positive rational numbers α_i and elementary atoms U_i satisfying (14). But now $S = U_1^{\alpha_1} S' = U_1^{\alpha_1} U_2^{\alpha_2} \cdot \ldots \cdot U_\ell^{\alpha_\ell}$ with (14) holding for $j \in [2, \ell]$. Since $g_1 \in \mathrm{supp}(U_1)$ and $g_1 \notin \mathrm{supp}(S') = \mathrm{supp}(U_2 \cdot \ldots \cdot U_\ell)$, we see that (14) holds when $j = 1$ and thus (14) is established. This completes the proof of the first part of the theorem. It remains to prove the upper bound for $\mathsf{D}(G_0)$.

Let $U \in \mathcal{A}(G_0)$ be an atom. We must show that $|U|$ is at most the bound given at the end of Theorem 3.7. If $|U| \leq 2$, this is clearly the case and so we may assume that $|U| \geq 3$. In this case, we may assume that $\mathrm{supp}(U) \cap \mathrm{supp}(-U) = \emptyset$. Let

$$U = U_1^{\alpha_1} \cdot \ldots \cdot U_\ell^{\alpha_\ell}$$

be the rational product decomposition of U given by the first part of the theorem. In particular, each $U_i \in \mathcal{A}(G_0)$ is an elementary atom and each $\alpha_i \in \mathbb{Q}_{>0}$ is a

positive rational number. We note that the corresponding sequence R is trivial since $\mathrm{supp}(U) \cap \mathrm{supp}(-U) = \emptyset$. If $\alpha_i > 1$ for some $i \in [1, \ell]$, then $U_i \mid U$ is a proper nontrivial zero-sum subsequence, contradicting that $U \in \mathcal{A}(G_0)$ is an atom. Thus $\alpha_i \leq 1$ for all $i \in [1, \ell]$. Now

$$|U| = \sum_{i=1}^{\ell} \alpha_i |U_i| \leq \sum_{i=1}^{\ell} |U_i| \leq \ell \mathsf{D}^{\mathrm{elm}}(G_0),$$

and, noting that $\mathrm{supp}(U) = \frac{1}{2}|\mathrm{supp}^+(U)|$ (since $\mathrm{supp}(U) \cap \mathrm{supp}(-U) = \emptyset$), the desired bound for $|U|$ follows from the bound for ℓ given in the first part of the theorem. □

We now consider which subsets $X \subset G_0 \subset \mathbb{Z}^r$ can be attained as the support of an elementary zero-sum sequence. A related question asks which subsets $X \subset G_0 \cup -G_0 \subset \mathbb{Z}^r$ can be attained as the signed support of an elementary zero-sum sequence. Given any $U \in \mathcal{B}_{\mathrm{rat}}(G_0)$, we know that $U^n \in \mathcal{B}(G_0)$ for some $n \geq 1$. Applying Lemma 3.5, we see that $X = \mathrm{supp}^+(U)$ for some elementary $U \in \mathcal{B}(G_0)$ is equivalent to $X = \mathrm{supp}^+(U)$ for some elementary $U \in \mathcal{B}_{\mathrm{rat}}(G_0)$ which in turn is equivalent to $X = \mathrm{supp}^+(U)$ for some elementary $U \in \mathcal{A}(G_0)$. The same is true when considering $X = \mathrm{supp}(U)$ for an elementary zero-sum sequence U over G_0. Of course, if $X = \mathrm{supp}^+(U)$ for a zero-sum sequence U, then $X \subset G_0 \cup -G_0$ is symmetric and so $X = Y \cup -Y$ for some $Y \subset G_0$ with $Y \cap -Y = \emptyset$. The following lemma classifies the possibilities for X.

Lemma 3.8. *Let $r \geq 1$, let $G_0 \subset \mathbb{Z}^r$ be a nonempty subset, and let $X \subset G_0 \cup -G_0$ be a subset with $X = -X$. Then condition (a) holds if and only if conditions (b) and (c) both hold. If $G_0 = -G_0$, then (a) and (b) are equivalent.*

(a) *$X = \mathrm{supp}^+(U)$ for some elementary zero-sum sequence $U \in \mathcal{B}(G_0)$.*
(b) *The elements of $X \cap G_0^+$ are linearly dependent over $\mathbb{Q}$, but any proper subset of $X \cap G_0^+$ is linearly independent over $\mathbb{Q}$.*
(c) *There exists a nontrivial zero-sum sequence $S \in \mathcal{B}(G_0)$ with $\emptyset \neq \mathrm{supp}^+(S) \subset X$.*

In particular, if $U \in \mathcal{B}(G_0)$ is an elementary zero-sum sequence such that $\mathrm{supp}(S) \cap \mathrm{supp}(-S) = \emptyset$, then $|\mathrm{supp}(U)| \leq r + 1$.

Proof. We set $Y = X \cap G_0^+$. Suppose that (a) holds. Then (c) holds and for each $g \in Y$ there exists a nonzero $\alpha_g \in \mathbb{Q} \backslash \{0\}$ such that

$$\sum_{g \in Y} \alpha_g g = 0 \quad \text{with } \alpha_g > 0 \text{ whenever } -g \notin G_0. \tag{15}$$

Since each α_g is nonzero, the elements of Y are linearly dependent.

Now, if we assume that condition (b) does not hold, then there must be some nonempty, proper subset $Z \subsetneq Y$ such that the elements of Z are linearly dependent. But then for each $g \in Z$ there exists $\beta_g \in \mathbb{Q}$ such that

$$\sum_{g \in Z} \beta_g g = 0 \tag{16}$$

with not all β_g zero. Set $\beta_g = 0$ for any $g \in Y \backslash Z$. Suppose first that there is some $g \in Y$ with $\alpha_g > 0$ and $\beta_g \neq 0$. Multiplying the β_g by -1 if need be, we may assume $\alpha_g > 0$ and $\beta_g > 0$. In this case, let $\gamma = \min\{\alpha_g/\beta_g : \alpha_g > 0, \beta_g > 0, g \in Y\} > 0$. If $\beta_g = 0$ whenever $\alpha_g > 0$, we set $\gamma = \min\{\alpha_g/\beta_g : \alpha_g < 0, \beta_g < 0, g \in Y\} > 0$. Multiplying (16) by γ, we obtain

$$\sum_{g \in Z} \gamma \beta_g g = 0. \tag{17}$$

Moreover, by the definition of γ we see that $\gamma \beta_g \leq \alpha_g$ whenever $\alpha_g > 0$ and $\beta_g > 0$. Thus, if we subtract (17) from (15), the resulting coefficient $\alpha_g - \gamma \beta_g$ will be nonnegative whenever $\alpha_g > 0$. As a result, since $\alpha_g > 0$ whenever $-g \notin G_0$ (by (15)), we see that

$$\sum_{g \in Z} (\alpha_g - \gamma \beta_g) g = 0 \quad \text{with } \alpha_g - \gamma \beta_g \geq 0 \text{ whenever } -g \notin G_0. \tag{18}$$

Furthermore, for an element $g_1 \in Y$ attaining the minimum in the definition of γ, we see that the coefficient $\alpha_{g_1} - \gamma \beta_{g_1}$ of g_1 in (18) is zero, while not all coefficients in (18) are zero since each α_g is nonzero and since at least one β_g is zero (Z is a proper subset of Y). Thus the $\mathbb{Q}$-linear relation given in (18) corresponds to a nontrivial rational zero-sum sequence $V \in \mathcal{B}_{\mathrm{rat}}(G_0)$ whose support is strictly contained in the support of U, contradicting that $U \in \mathcal{B}(G_0)$ is an elementary zero-sum sequence. We can then conclude that all proper subsets of Y are linearly independent, as desired.

Now suppose that (b) holds and that either $G_0 = -G_0$ or (c) holds. Clearly, there cannot be any zero-sum sequence U over $G_0 \subset \mathbb{Z}^r \subset \mathbb{Q}^r$ with $\mathrm{supp}(U)$ consisting of linearly independent elements over $\mathbb{Q}$. Thus, in order to show that (a) holds, it suffices to show that there exists some $U \in \mathcal{B}_{\mathrm{rat}}(G_0)$ with $\mathrm{supp}^+(U) = X$.

If (c) holds, then a nontrivial $U \in \mathcal{B}_{\mathrm{rat}}(G_0)$ exists with $\emptyset \neq \mathrm{supp}^+(U) \subset X$. However, since $\mathrm{supp}^+(U) \cap G^+$ can only be linearly independent if $\mathrm{supp}^+(U)$ is empty, we have $\mathrm{supp}^+(U) = X$ as desired.

Next assume $G_0 = -G_0$. We need to show that there exists some $U \in \mathcal{B}_{\mathrm{rat}}(G_0)$ with $\mathrm{supp}^+(U) = X$. This is equivalent to showing that there exists some nonzero $\alpha_1, \ldots, \alpha_\ell \in \mathbb{Q} \backslash \{0\}$ with $\sum_{i=1}^{\ell} \alpha_i g_i = 0$. As the elements of Y are linearly dependent, there exist $\alpha_1, \ldots, \alpha_\ell \in \mathbb{Q}$ not all zero with $\sum_{i=1}^{\ell} \alpha_i g_i = 0$. Thus, if $\alpha_i = 0$ for some $i \in [1, \ell]$, this would give a dependence relation on the elements of $Y \backslash \{g_i\}$, contradicting the hypothesis that every proper subset of Y is linearly independent. Consequently, $\alpha_i \neq 0$ for all $i \in [1, \ell]$, and (a) follows as noted earlier.

If $U \in \mathcal{B}(G_0)$ is an elementary zero-sum sequence such that $\mathrm{supp}(U) \cap \mathrm{supp}(-U) = \emptyset$, then the first part of the theorem implies that $\mathrm{supp}(U) \backslash \{g\} \subset \mathbb{Q}^r$ is a set of linearly independent vectors for any $g \in \mathrm{supp}(U)$. Since any subset of vectors of size $r+1$ must be linearly dependent in $\mathbb{Q}^r$, the desired bound $|\mathrm{supp}(U)| \leq r+1$ follows. □

Lemma 3.9. *Let $r \geq 1$ and let $G_0 \subset \mathbb{Z}^r$ be a nonempty subset. Then $\mathsf{D}(G_0) \geq 3$ if and only if there exists an elementary atom $U \in \mathcal{A}(G_0)$. If this is the case, then $\mathsf{D}^{\mathrm{elm}}(G_0) \geq 3$ as well.*

Proof. Since any elementary atom $U \in \mathcal{A}(G_0)$ must satisfy $|U| \geq 3$, one direction is clear. Suppose now that $\mathsf{D}(G_0) \geq 3$ and let $V \in \mathcal{A}(G_0)$ be an atom with $|V| \geq 3$. Then $\mathrm{supp}^+(V)$ is nonempty, in which case Lemma 3.6 completes the proof. □

Recall that, given any $m \times n$ integer matrix M with $m \leq n$, we can perform elementary row and column operations on M (swapping rows/columns, multiplying a row/column by ± 1, or adding an integer multiple of a row/column to another row/column) to obtain a diagonal integer matrix $D = (d_{i,j})_{i,j}$ with $d_{1,1}|\cdots|d_{m,m}$ and $d_{i,i} \geq 0$ for all $i \in [1,m]$. The matrix D is unique and is known as the Smith normal form of the matrix M and the $d_{i,i}$ are called the *elementary divisors* of M. If $g_1, \ldots, g_n \in \mathbb{Z}^m$ are the columns of M, then $\mathbb{Z}^m / \langle g_1, \ldots, g_n \rangle \cong \mathbb{Z}/d_{1,1}\mathbb{Z} \oplus \cdots \oplus \mathbb{Z}/d_{m,m}\mathbb{Z}$. Thus, when M has full rank (whence $d_{m,m} \neq 0$), we have $d_{m,m} = \exp(\mathbb{Z}^m / \langle g_1, \ldots, g_n \rangle)$. It is easily checked (and well-known) that for $j \in [1,m]$

$$\gcd\{\det(M') : M' \text{ is a } j \times j \text{ submatrix of } M\}$$

remains invariant under elementary row and column operations and is thus equal to $d_{1,1} \cdots d_{j,j}$. In particular, if $m = n$, so that M is a square matrix,

$$d_{m,m} = \frac{|\det(M)|}{\gcd\{|\det(M')| : M' \text{ is a } (m-1) \times (m-1) \text{ submatrix of } M\}}. \qquad (19)$$

These results can be found in many standard textbooks dealing with linear algebra over $\mathbb{Z}$.

We now turn our attention to finding bounds for $\mathsf{D}(G_0)$ where $G_0 \subset \mathbb{Z}^r$. Let M be a $r \times |G_0^+|$ matrix whose columns are the vectors $g \in G_0^+ \subset \mathbb{Z}^r$. Using lattice theory and results from the geometry of numbers, Diaconis, Graham, and Sturmfels [5] showed that

$$\mathsf{D}(G_0) \leq (2r)^r(r+1)^{r+1} \max\{|\det(M')| : M' \text{ is a } r \times r \text{ submatrix of } M\}$$

when G_0 is finite with full rank $\mathsf{r}(\langle G_0 \rangle) = r$. However, when $|G_0|$ is not terribly large, Theorem 3.7 can be used to obtain tighter bounds. To do so, we need to be able to bound $\mathsf{D}^{\mathrm{elm}}(G_0)$ and to do this we consider an argument of Sturmfels [52, Chap. 4] which, when combined with additional results, allows us to give a linear algebraic description of $\mathsf{D}^{\mathrm{elm}}(G_0)$. In order to state the next theorem we first need

the following definition. For a collection of $r + 1$ integer vectors $g_1, \ldots, g_{r+1} \in \mathbb{Z}^r$, we define

$$\Delta(g_1, \ldots, g_{r+1}) = \frac{\sum_{i=1}^{r+1} |\det(g_1, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1})|}{\gcd\{|\det(g_1, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1})| : i \in [1, r+1]\}} \in \mathbb{N}_0,$$

where $\Delta(g_1, \ldots, g_{r+1}) = 0$ if $\mathsf{r}(\langle g_1, \ldots, g_{r+1} \rangle) < r$.

Theorem 3.10. *Let $r \geq 1$ and let $G_0 \subset \mathbb{Z}^r$ be a nonempty subset with $\mathsf{r}(\langle G_0 \rangle) = r$ and $\mathsf{D}(G_0) \geq 3$. Then*

$$\mathsf{D}^{\mathrm{elm}}(G_0) = \sup\{\Delta(g_1, \ldots, g_{r+1}) : g_1, \ldots, g_{r+1} \in G_0 \text{ and } \mathsf{D}(\{g_1, \ldots, g_{r+1}\}) \geq 3\}$$

$$\leq \sup\{\Delta(g_1, \ldots, g_{r+1}) : g_1, \ldots, g_{r+1} \in G_0\}.$$

Moreover, if $G_0 = -G_0$, then equality holds.

Proof. By Lemma 3.9, the hypothesis $\mathsf{D}(G_0) \geq 3$ is equivalent to the existence of an elementary atom $U \in \mathcal{A}(G_0)$. Since any such elementary atom U satisfies $|U| \geq 3$, we conclude that

$$\mathsf{D}^{\mathrm{elm}}(G_0) \geq 3. \tag{20}$$

Let $U \in \mathcal{A}(G_0)$ be an arbitrary elementary atom and let $X = \mathrm{supp}(U)$. Since U is an elementary atom, $\mathrm{supp}(U) \cap \mathrm{supp}(-U) = \emptyset$. By Lemma 3.8, X is linearly dependent over $\mathbb{Q}$, but any proper subset of X is linearly independent. In particular, $|X| = x + 1 \geq 2$ with

$$1 \leq \mathsf{r}(\langle X \rangle) = x \leq \mathsf{r}(\langle G_0 \rangle) = r.$$

Thus, if $x < r$, then we can find a subset $X' \subset G_0 \backslash (X \cup -X)$ such that $|X'| = r - x$ and $\mathsf{r}(\langle X \cup X' \rangle) = r$. Let $Y = X \cup X'$. Note that $\mathrm{supp}(U) \subset Y \subset G_0$ with $|Y| = r + 1$, and that, for each $g \in G_0, g$ and $-g$ are not both contained in Y. Let $Y = \{g_1, \ldots, g_{r+1}\}$ where $g_1, \ldots, g_{x+1}$ are the elements from X.

Let M be the $r \times (r + 1)$ matrix whose columns are the vectors $g_i \in Y \subset G_0$. Then the vector of integer multiplicities $\mathbf{x} = (x_i)_{i \in [1, r+1]} \in \mathbb{Z}^{|Y|} = \mathbb{Z}^{r+1}$ corresponds to a zero-sum subsequence $S = \prod_{i=1}^{r+1} g_i^{x_i} \in \mathcal{B}(G_0)$ with $\mathrm{supp}(S) \subset Y$ (with $\mathrm{supp}(S) \cap \mathrm{supp}(-S) = \emptyset$) and $\mathsf{v}_{g_i}(S) = x_i$ precisely when $\mathbf{x}$ is in the kernel of the matrix M and $x_i \geq 0$ for all $i \in [1, r+1]$. Also, the vector $\mathbf{x} \in \mathbb{Z}^{r+1}$ corresponds to a zero-sum subsequence $S = \prod_{i \in I^+} g_i^{x_i} \prod_{i \in I^-} (-g_i)^{-x_i} \in \mathcal{B}(G_0 \cup -G_0)$ with $\mathrm{supp}^+(S) \subset Y \cup -Y$ precisely when $\mathbf{x}$ is in the kernel of the matrix M, $I^+ \subset [1, r+1]$ denotes the subset of indices $i \in [1, r+1]$ with $x_i > 0$, and $I^- \subset [1, r+1]$ denotes the subset of indices $i \in [1, r+1]$ with $x_i < 0$. In the latter case, we have $S \in \mathcal{B}(G_0)$ precisely when $-Y^- \subset G_0$ for $Y^- = \{y_i \in Y : i \in I^-\}$.

Consider the vector $\mathbf{x} = \frac{1}{\delta}(x_i)_{i \in [1, r+1]} \in \mathbb{Z}^{r+1}$ given by

$$x_i = (-1)^i \det(g_1, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}) \quad \text{for } i \in [1, r+1],$$

where

$$\delta = \gcd\{\det(g_1, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}) : i \in [1, r+1]\}.$$

Since $\mathsf{r}(\langle Y \rangle) = r$, the x_i cannot all be zero, and thus $\delta > 0$ is also nonzero. Since each $g_i \in \mathbb{Z}^r$ is an integer-valued vector, it is clear from the above definition that $\mathbf{x} \in \mathbb{Z}^{r+1}$ with

$$\gcd\left\{ \frac{x_i}{\delta} : i \in [1, r+1] \right\} = 1. \tag{21}$$

Moreover, since $X = \{g_1, \ldots, g_{x+1}\}$ is linearly dependent, $x_i = 0$ for all $i \in [x+2, r+1]$.

We now show that $\mathbf{x}$ is in the kernel of the matrix M whose columns are the vectors $g_i \in \mathbb{Z}^r$. Let $g_i = (g_{i,j})_{j \in [1,r]}$ with $g_{i,j}$ the jth entry of the column vector g_i. With $j \in [1, r]$ arbitrary, the jth entry of $M\mathbf{x} \in \mathbb{Z}^r$ is

$$\frac{1}{\delta} \sum_{i=1}^{r+1} g_{i,j} x_i = \frac{1}{\delta} \sum_{i=1}^{r+1} (-1)^i g_{i,j} \det(g_1, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}). \tag{22}$$

However, by the cofactor expansion formula for the determinant of a matrix, the right-hand side of (22) is equal (up to sign) to the product of $\frac{1}{\delta}$ and the determinant of the $(r+1) \times (r+1)$ matrix M' formed from the $r \times (r+1)$ matrix M by repeating the jth row $(g_{i,j})_{i \in [1,r]}$ of M and then computing the cofactor expansion about this duplicate row. As M' has two duplicate rows, its determinant is zero and thus (22) is zero if $j \in [1, r]$. Since the jth entry of $M\mathbf{x} \in \mathbb{Z}^r$ is equal to (22), this shows that every coordinate of $M\mathbf{x}$ is zero. Hence $M\mathbf{x} = 0$ and $\mathbf{x}$ is in the kernel of M as claimed.

Since $\mathbf{x} = \frac{1}{\delta}(x_i)_{i \in [1, r+1]} \in \mathbb{Z}^{r+1}$ is an integer-valued vector in the kernel of M, it follows (as was noted earlier in the proof) that $\mathbf{x}$ corresponds to a nontrivial (not all x_i are zero) zero-sum sequence $S \in \mathcal{B}(G_0 \cup -G_0)$ with $\mathrm{supp}^+(S) \subset X \cup -X$ ($x_i = 0$ for $i \in [x+2, r+1]$ and $X = \{g_1, \ldots, g_{x+1}\}$). Moreover, $\mathrm{supp}(S) \cap \mathrm{supp}(-S) = \emptyset$. Consequently, as $X \cup -X = \mathrm{supp}^+(U)$ with $U \in \mathcal{A}(G_0)$ elementary, $\mathrm{supp}^+(S) = X \cup -X$. From (21) we see that $S \neq T^\ell$ for any $T \in \mathcal{B}(G_0 \cup -G_0)$ and $\ell \geq 2$. Since any proper subset of $\mathrm{supp}(U)$ is linearly independent, it follows that there is no nontrivial zero-sum sequence $V \in \mathcal{B}(G_0 \cup -G_0)$ with $\emptyset \neq \mathrm{supp}^+(V) \subsetneq \mathrm{supp}^+(U)$. Thus U is an elementary atom not just over G_0, but also over $G_0 \cup -G_0$. From Lemma 3.5, U must be the unique (up to sign) elementary atom over $G_0 \cup -G_0$ with signed support $X \cup -X$, and all other elementary zero-sum sequences T over $G_0 \cup -G_0$ with $\mathrm{supp}^+(T) = X \cup -X$ (for which $\mathrm{supp}(T) \cap \mathrm{supp}(-T) = \emptyset$) must be a power of either U or $-U$. Applying this conclusion to $T = S$, we find that either $S = U^\ell$ or $-S = U^\ell$ for some $\ell \geq 1$. By swapping the sign of each x_i in the definition of $\mathbf{x} = (x_i)_{i \in [1,r]}$ (thus replacing $\mathbf{x}$ by $-\mathbf{x}$) if need be, we may without loss of generality assume the former: $S = U^\ell$ and, in particular, $S \in \mathcal{B}(G_0)$. Since $S \neq T^\ell$ for any $T \in \mathcal{B}(G_0 \cup -G_0)$ and $\ell \geq 2$ as observed above, it follows that $\ell = 1$ and $S = U$.

Now

$$|U| = |S| = \frac{1}{\delta} \sum_{i=1}^{r} |x_i| = \frac{1}{\delta} \sum_{i=1}^{r+1} |\det(g_1, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1})|$$

$$= \Delta(g_1, \ldots, g_{r+1}).$$

Since $U \in \mathcal{A}(G_0)$ was an arbitrary elementary atom with $\operatorname{supp}(U) \subset \{g_1, \ldots, g_{r+1}\}$, we have

$$\mathsf{D}^{\mathsf{elm}}(G_0) \leq \sup\{\Delta(g_1, \ldots, g_{r+1}) : g_1, \ldots, g_{r+1} \in G_0 \quad \text{and}$$
$$\mathsf{D}(\{g_1, \ldots, g_{r+1}\}) \geq 3\} \tag{23}$$
$$\leq \sup\{\Delta(g_1, \ldots, g_{r+1}) : g_1, \ldots, g_{r+1} \in G_0\}. \tag{24}$$

Let $g_1, \ldots, g_{r+1} \in G_0$ be vectors with $\mathsf{D}(\{g_1, \ldots, g_{r+1}\}) \geq 3$ and such that $\Delta(g_1, \ldots, g_{r+1}) > 0$. Such vectors exist by (20) and (23). The condition $\Delta(g_1, \ldots, g_{r+1}) > 0$ implies that $\mathsf{r}(\langle g_1, \ldots, g_{r+1}\rangle) = r$. Since $\mathsf{D}(\{g_1, \ldots, g_{r+1}\}) \geq 3$ and by Lemma 3.9, there exists an elementary atom $U \in \mathcal{A}(\{g_1, \ldots, g_{r+1}\})$, such that, without loss of generality, $\operatorname{supp}(U) = \{g_1, \ldots, g_{x+1}\}$, where $x \leq r$. Now, repeating the above arguments using this particular elementary atom U, we find that $\mathsf{D}^{\mathsf{elm}}(G_0) \geq |U| = \Delta(g_1, \ldots, g_{r+1})$. Taking the supremum over $\Delta(g_1, \ldots, g_{r+1})$ for all choices of $g_1, \ldots, g_{r+1} \in G_0$ with $\mathsf{D}(\{g_1, \ldots, g_{r+1}\}) \geq 3$ and $\Delta(g_1, \ldots, g_{r+1}) > 0$, we see that equality holds in (23).

Next suppose that $-G_0 = G_0$. To complete the proof we need to show that equality holds in (24). We may assume $\mathsf{D}^{\mathsf{elm}}(G_0) < \infty$ as the claim is trivially true otherwise. Let $g_1, \ldots, g_{r+1} \in G_0$ be vectors that obtain the maximum in (23). If $\mathsf{r}(\langle g_1, \ldots, g_{r+1}\rangle) < r$, then $\Delta(g_1, \ldots, g_{r+1}) = 0$ and, by the lower bound $\mathsf{D}^{\mathsf{elm}}(G_0) \geq 3$ from (20), the $g_1, \ldots, g_{r+1} \in G_0$ cannot maximize (24). Thus we may assume that $\mathsf{r}(\langle g_1, \ldots, g_{r+1}\rangle) = r$. Note that replacing any g_i with $-g_i$ does not alter the value of $\Delta(g_1, \ldots, g_{r+1})$ and that the hypothesis $G_0 = -G_0$ ensures that $-g_i \in G_0$. Thus, to show equality in (24), it suffices to show that $\mathsf{D}(\{\epsilon_1 g_1, \ldots, \epsilon_{r+1} g_{r+1}\}) \geq 3$ for some choice of $\epsilon_i \in \{1, -1\}$ which, by Lemma 3.9, is equivalent to the existence of an elementary atom $U \in \mathcal{A}(\{\epsilon_1 g_1, \ldots, \epsilon_{r+1} g_{r+1}\})$ for some choice of $\epsilon_i \in \{1, -1\}$ which in turn is equivalent to the existence of an elementary atom $U \in \mathcal{A}(G_0)$ with $\operatorname{supp}^+(U) \subset \{g_1, \ldots, g_{r+1}\} \cup -\{g_1, \ldots, g_{r+1}\}$. To show the later we will use Lemma 3.8.

Let $Y \subset \{g_1, \ldots, g_{r+1}\}$ be those $g_i \in \{g_1, \ldots, g_{r+1}\}$ such that $\det(g_1, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}) \neq 0$. It follows that $\{g_1, \ldots, g_{r+1}\} \backslash \{g_i\}$ is linearly independent for any $g_i \in Y$. Thus, any subset of $\{g_1, \ldots, g_{r+1}\} \backslash \{g_i\}$, including $Y \backslash \{g_i\}$, must also be linearly independent. This shows that all proper subsets of Y are linearly independent.

Suppose the elements of Y are linearly independent. Then clearly $|Y| \leq r$ and thus Y must be a proper subset of $\{g_1, \ldots, g_{r+1}\}$. However, since $\mathsf{r}(\langle g_1, \ldots, g_{r+1}\rangle) = r$, it follows that we can complete Y to some full rank set $Y' \subset \{g_1, \ldots, g_{r+1}\}$ with

$Y \subset Y'$. Then $\det(g_1, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}) \neq 0$ for the unique $g_i \in \{g_1, \ldots, g_{r+1}\} \backslash Y'$. By the definition of Y, this forces $g_i \in Y$, contradicting that $g_i \notin Y'$ with $Y \subset Y'$. Thus we conclude that the elements of Y are linearly dependent.

In view of the conclusions of the previous two paragraphs, along with the hypothesis $G_0 = -G_0$, we can now apply Lemma 3.8 and conclude that there exists an elementary zero-sum sequence $S \in \mathcal{B}(G_0)$ with $\mathrm{supp}^+(S) = Y \subset \{g_1, \ldots, g_{r+1}\} \cup -\{g_1, \ldots, g_{r+1}\}$. By Lemma 3.5, this ensures that there exists an elementary atom $U \in \mathcal{A}(G_0)$ with $\mathrm{supp}^+(U) = Y \subset \{g_1, \ldots, g_{r+1}\} \cup -\{g_1, \ldots, g_{r+1}\}$ which completes the proof. $\qquad\square$

Corollary 3.11. *Let $r \in \mathbb{N}$, $(\mathsf{e}_1, \ldots, \mathsf{e}_{r+1})$ denote the standard basis of $\mathbb{Z}^{r+1}$, and $G_0 \subset G = \langle \mathsf{e}_1, \ldots, \mathsf{e}_r \rangle$ be a nonempty subset with $\mathsf{r}(\langle G_0 \rangle) = r$ and $\mathsf{D}(G_0) \geq 3$. Furthermore, let $\widetilde{G}_0, \widetilde{G}_1 \subset \mathbb{Z}^{r+1}$ be the subsets given by*

$$\widetilde{G}_0 = \{g + \mathsf{e}_{r+1} : g \in G_0\} \cup G_0 \quad \text{and}$$

$$\widetilde{G}_1 = \{g + \mathsf{e}_{r+1} : g \in G_0\} \cup \{g - \mathsf{e}_{r+1} : g \in G_0\},$$

and let $\mathcal{M}(\widetilde{G}_i)$ be all those nonsingular $(r+1) \times (r+1)$ matrices with columns $\tilde{g}_1, \ldots, \tilde{g}_{r+1} \in \widetilde{G}_i \subset \mathbb{Z}^{r+1}$, for $i \in [0,1]$. Also, let $d_{r+1}(\tilde{M})$ denote the largest elementary divisor of the matrix $\tilde{M}$.

(1) $\mathsf{D}^{\mathsf{elm}}(G_0) \leq 2 \sup\{d_{r+1}(\tilde{M}) : \tilde{M} \in \mathcal{M}(\widetilde{G}_0)\} \leq 2 \sup\{|\det(\tilde{M})| : \tilde{M} \in \mathcal{M}(\widetilde{G}_0)\}$
    and
(2) $\mathsf{D}^{\mathsf{elm}}(G_0) \leq \sup\{d_{r+1}(\tilde{M}) : \tilde{M} \in \mathcal{M}(\widetilde{G}_1)\} \leq \sup\{|\det(\tilde{M})| : \tilde{M} \in \mathcal{M}(\widetilde{G}_1)\}$.

Proof. By Lemma 3.9 and Theorem 3.10 we know that

$$3 \leq \mathsf{D}^{\mathsf{elm}}(G_0) \leq \sup\{\Delta(g_1, \ldots, g_{r+1}) : g_1, \ldots, g_{r+1} \in G_0\}.$$

It follows that there exist $g_1, \ldots, g_{r+1} \in G_0 \subset \mathbb{Z}^r$ with $\Delta(g_1, \ldots, g_{r+1}) > 0$ and that the supremum on the right-hand side is necessarily obtained for such a choice of $g_1, \ldots, g_{r+1}$. Let $g_1, \ldots, g_{r+1} \in G_0$ be such that $\Delta(g_1, \ldots, g_{r+1}) > 0$. Then $\mathsf{r}(\langle g_1, \ldots, g_{r+1} \rangle) = r$.

For each $i \in [1, r+1]$, let $\tilde{g}_i = g_i \pm \mathsf{e}_{r+1}$, where an appropriate choice for the sign of e_{r+1} will be determined shortly, and let $\tilde{M} \in \mathcal{M}(\widetilde{G}_1)$ be the $(r+1) \times (r+1)$ integer matrix with columns $\tilde{g}_1, \ldots, \tilde{g}_{r+1} \in \tilde{G}_1 \subset \mathbb{Z}^{r+1}$. Let

$$\delta = \gcd\{|\det(g_1, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1})| : i \in [1, r+1]\} \quad \text{and}$$

$$\tilde{\delta} = \gcd\{|\det(M')| : M' \text{ is a } r \times r \text{ sub-matrix of } \tilde{M}\}.$$

Then $\tilde{\delta} \mid \delta$. By using the cofactor expansion formula for the determinant and expanding along the final row of $\tilde{M}$, we see that by choosing the signs in each $\tilde{g}_i$ appropriately, $\frac{1}{\delta} \det(\tilde{M}) = \Delta(g_1, \ldots, g_{r+1})$. Now, since $\tilde{\delta} \mid \delta$ and by (19), we have

$$\Delta(g_1, \ldots, g_{r+1}) = \frac{1}{\delta} \det(\tilde{M}) \leq \frac{1}{\tilde{\delta}} \det(\tilde{M}) = d_{r+1}(\tilde{M}).$$

Since the choice of $g_1, \ldots, g_{r+1} \in G_0$ was arbitrary, subject to the restriction $\Delta(g_1, \ldots, g_{r+1}) > 0$, the bound from Theorem 3.10 establishes (2).

For each $i \in [1, r+1]$, redefine $\tilde{g}_i$ as either $\tilde{g}_i = g_i + e_{r+1}$ or $\tilde{g}_i = g_i$, where the choice of coefficient for e_{r+1} will be determined shortly, and let $\tilde{M} \in \mathcal{M}(\widetilde{G}_0)$ be the $(r+1) \times (r+1)$ integer matrix with columns $\tilde{g}_1, \ldots, \tilde{g}_{r+1} \in \tilde{G}_0 \subset \mathbb{Z}^{r+1}$. Let

$$\delta = \gcd\{|\det(g_1, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1})| : i \in [1, r+1]\} \quad \text{and}$$

$$\tilde{\delta} = \gcd\{|\det(M')| : M' \text{ is a } r \times r \text{ sub-matrix of } \tilde{M}\}.$$

Then $\tilde{\delta} \,|\, \delta$. By again using the cofactor expansion formula for the determinant and expanding along the final row of $\tilde{M}$, we see that by choosing the coefficients for e_{r+1} in each $\tilde{g}_i$ appropriately, we can achieve

$$\det(\tilde{M}) = \sum_{i \in I} |\det(g_1, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1})|$$

for some subset $I \subset [1, r]$. By using the exact opposite choices for the $\tilde{g}_i$, we can instead achieve

$$\det(\tilde{M}) = - \sum_{i \in [1, r+1] \setminus I} |\det(g_1, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1})|.$$

Between these two options, take the one where $|\det(\tilde{M})|$ is larger. Then

$$2|\det(\tilde{M})| \geq \sum_{i=1}^{r+1} |\det(g_1, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1})|,$$

in which case (19) and $\tilde{\delta} \,|\, \delta$ give

$$\Delta(g_1, \ldots, g_{r+1}) \leq \frac{2}{\delta} |\det(\tilde{M})| \leq \frac{2}{\tilde{\delta}} |\det(\tilde{M})| = 2d_{r+1}(\tilde{M}),$$

establishing (1). $\qquad \square$

The special case $G_0 = G_r^+ \cup -G_r^+$. Our goal for the remainder of this section is to apply the machinery above to the case when

$$G_0 = G_r^+ \cup -G_r^+,$$

where G_r^+ denotes the nonzero vertices of the r-dimensional hypercube as defined in Sec. 2. Specifically, we wish to obtain upper and lower bounds for $\mathsf{D}(G_r^+ \cup -G_r^+)$ which we will apply to the study of invariants of monoids of modules in Sec. 4.

If $r = 1$, then clearly $\mathsf{D}(G_0) = 2$. Thus we suppose for the remainder of this section that $r > 1$ in which case $\mathsf{D}(G_0) \geq 3$. Note that $\Delta(g_1, \ldots, g_{r+1})$ is unaffected if any g_i is replaced by $-g_i$. Thus we need only consider matrices with columns from G_r^+ when applying Theorem 3.10 and Corollary 3.11. By Corollary 3.11 we see that $\mathsf{D}^{\mathsf{elm}}(G_0)$ is bounded from above by twice the maximal determinant of a $(r+1) \times (r+1)$ $(0, 1)$-matrix. It is well-known that this value is $\frac{1}{2^{r+1}}$ times the maximal value of a $(r+2) \times (r+2)$ $(1, -1)$-matrix and this value, in turn, is equal

to the maximal value of a $(r+2) \times (r+2)$ $(-1, 0, 1)$-matrix. Moreover, this maximal value is bounded by Hadamard's bound and, together with Theorem 3.7, we obtain the following upper bound

$$\mathsf{D}(G_0) \leq (|G_0^+| - r)\mathsf{D}^{\text{elm}}(G_0) \leq \frac{(2^r - r - 1)}{2^r}(r + 2)^{\frac{r+2}{2}} < (r + 2)^{\frac{r+2}{2}}.$$

We now construct a lower bound for $\mathsf{D}(G_0)$. Suppose we have an $(r + 1) \times (r + 1)$ $(0, 1)$-matrix M whose determinant is a large prime p. Then we must have $d_{r+1} = p$ and $d_r = 1$ where the $d_1 | \cdots | d_{r+1}$ are the elementary divisors of M. Therefore, by (19), we see that the greatest common divisor of the determinants of $r \times r$ sub-matrices of M must be 1. On the other hand, if we delete any row of M, we obtain, by applying the cofactor expansion formula to $\det(M)$ and expanding along the deleted, row $r + 1$ vectors $g_1, \ldots, g_{r+1} \in G_0$ where $\gcd\{|\det(g_1, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1})| : i \in [1, r + 1]\}$ divides $\det(M) = p$. As p is prime, we can conclude that $\gcd\{\det(g_1, \ldots, g_{i-1}, g_{i+1}, \ldots, g_{r+1}) : i \in [1, r + 1]\}$ is either 1 or p. Moreover, since the greatest common divisor of the determinants of $r \times r$ sub-matrices is 1, we see that we can achieve 1 rather than p by choosing an appropriate row of M to delete. Then, in this case, $\Delta(g_1, \ldots, g_{r+1}) = d_{r+1} = |\det(M)| = p$ and we obtain the lower bound

$$\max\{|\det(M)| : M \text{ is a } (r + 1) \times (r + 1) \ (0, 1)\text{-matrix with } |\det(M)| \text{ prime}\}$$

$$\leq \mathsf{D}^{\text{elm}}(G_0) \leq \mathsf{D}(G_0).$$

In general, the possible values of $|\det(M)|$ for an arbitrary $(0, 1)$-matrix are not known, though this question is of great interest to many researchers (see [46] for known results on the spectrum of the determinant). However, computational evidence obtained for small values of r has led many to observe that there is (at least for r small) a constant $C \cong \frac{1}{2}$ and a large consecutive interval $[0, C\,2^{-r-1}(r + 2)^{(r+2)/2}]$ of obtainable values for $|\det(M)|$, leading some to conjecture that the set of values of $\det(M)$ is dense in a interval whose length is a fraction of the maximal possible value (see [45]). For any interval of obtainable values $[0, n]$, Bertrand's postulate ensures that a prime of size at least $\frac{1}{2}n$ will occur in that interval. Thus, if the intuitions gathered from examining small values of r hold true for larger values in a very strong sense, we would expect a lower bound for $\mathsf{D}(G_0)$ of the form

$$\left(\frac{r + 2}{C}\right)^{(r+2)/2} \lesssim \mathsf{D}^{\text{elm}}(G_0) \leq \mathsf{D}(G_0),$$

where $C \geq 1$ is some constant. We note that such a lower bound would very nearly match the upper bound in order of magnitude.

Constructively, the best lower bounds we have been able to achieve involve the Fibonacci numbers. We now present a construction giving a lower bound on $\mathsf{D}(G_0)$ regardless of concerns about the possible values of $\det(M)$. We note that for small values of r, examples of $(0, 1)$-matrices with large prime determinant are known

and can thus be used to improve upon this bound. For $r \in \mathbb{N}_0$, we denote by F_r the rth *Fibonacci number*. That is, $\mathsf{F}_0 = 0, \mathsf{F}_1 = 1$, and $\mathsf{F}_r = \mathsf{F}_{r-1} + \mathsf{F}_{r-2}$ for all $r \geq 2$.

Proposition 3.12. *Let $r \in \mathbb{N}$, let $(\mathsf{e}_1, \ldots, \mathsf{e}_r)$ denote the standard basis of $\mathbb{Z}^r$, and set $H_r = \langle \mathsf{e}_1 + \cdots + \mathsf{e}_r \rangle$. Then there exists a sequence $S_r \in \mathcal{F}(G_r^+)$ with*

$$\sigma(S_r) \in H_r \quad and \quad \Sigma_{\leq |S_r|-1}(S_r) \cap H_r = \emptyset$$

such that

$$|\mathrm{supp}(S_r)| = r \quad with\ \mathrm{supp}(S_r)\ spanning\ \mathbb{Q}^r,$$
$$|S_r| = \mathsf{F}_{r+1} \quad and \quad \sigma(S_r) = \mathsf{F}_r \mathsf{e}_1 + \cdots + \mathsf{F}_r \mathsf{e}_r.$$

Proof. The sequence $S_1 = \mathsf{e}_1$ is easily seen to satisfy the conditions of the theorem. We proceed recursively to define S_r for $r \geq 2$, assuming that

$$S_{r-1} = g_1 \cdot \ldots \cdot g_{\mathsf{F}_r} \in \mathcal{F}(G_{r-1}^+)$$

has already been constructed so as to satisfy the conclusions of the theorem. Let

$$S_r' = S_{r-1} \mathsf{e}_r^{\mathsf{F}_{r-1}} \in \mathcal{F}(G_r^+)$$

and let

$$S_r = \varphi(S_r') \in \mathcal{F}(\mathbb{Z}^r),$$

where $\varphi \colon \mathbb{Z}^r \to \mathbb{Z}^r$ is the map defined by $x \mapsto -(x - (\mathsf{e}_1 + \cdots + \mathsf{e}_r))$ and extends to an affine linear isomorphism of $\mathbb{Q}^r$. The map φ acts on an element $x \in G_r^+$ simply by exchanging each 0 for a 1 and each 1 for a 0. Therefore, $\varphi(x) \in G_r^+$ for each $x \in G_r^+ \backslash \{\mathsf{e}_1 + \cdots + \mathsf{e}_r\}$. As a result, since $\mathsf{e}_1 + \cdots + \mathsf{e}_r \notin \mathrm{supp}(S_r')$ (as $\mathrm{supp}(S_{r-1}) \subset G_{r-1}^+$ and $r \geq 2$), we see that

$$S_r \in \mathcal{F}(G_r^+) \quad with\ |S_r| = |S_r'|\ and\ |\mathrm{supp}(S_r)| = |\mathrm{supp}(S_r')|.$$

Observe that

$$|S_r| = |S_r'| = |S_{r-1}| + \mathsf{F}_{r-1} = \mathsf{F}_r + \mathsf{F}_{r-1} = \mathsf{F}_{r+1}$$

and that

$$\sigma(S_r) = \sigma(\varphi(S_r')) = |S_r|\mathsf{e}_1 + \cdots + |S_r|\mathsf{e}_r - \sigma(S_r')$$
$$= \mathsf{F}_{r+1}\mathsf{e}_1 + \cdots + \mathsf{F}_{r+1}\mathsf{e}_r - (\mathsf{F}_{r-1}\mathsf{e}_1 + \cdots + \mathsf{F}_{r-1}\mathsf{e}_r)$$
$$= \mathsf{F}_r\mathsf{e}_1 + \cdots + \mathsf{F}_r\mathsf{e}_r \in H_r.$$

Moreover, since $\varphi \colon \mathbb{Q}^r \to \mathbb{Q}^r$ is an affine linear isomorphism, we have that

$$|\mathrm{supp}(S_r)| = |\mathrm{supp}(\varphi(S_r'))| = |\mathrm{supp}(S_r')| = |\mathrm{supp}(S_{r-1})| + 1 = r.$$

Since $\varphi \colon \mathbb{Q}^r \to \mathbb{Q}^r$ is an affine linear isomorphism, to show that $\mathrm{supp}(S_r) = \mathrm{supp}(\varphi(S_r'))$ spans $\mathbb{Q}^r$ it suffices to show that $\mathrm{supp}(S_r')$ spans $\mathbb{Q}^r$. But this is clear since $\mathrm{supp}(S_r') = \mathrm{supp}(S_{r-1}) \cup \{\mathsf{e}_r\}$ with $\mathrm{supp}(S_{r-1})$ spanning $\mathbb{Q}\mathsf{e}_1 + \cdots + \mathbb{Q}\mathsf{e}_{r-1}$ by hypothesis. It remains to show that $\Sigma_{\le |S_r|-1}(S_r) \cap H_r = \emptyset$. Since $\sigma(\varphi(T)) \in H_r$ if and only if $\sigma(T) \in H_r$ for any sequence $T \in \mathcal{F}(\mathbb{Z}^r)$, we see that in order to show $\Sigma_{\le |S_r|-1}(S_r) \cap H_r = \emptyset$, it suffices to show that $\Sigma_{\le |S_r'|-1}(S_r') \cap H_r = \emptyset$.

Suppose that $T \mid S_r$ is a nontrivial subsequence with $\sigma(T) \in H_r$, that is, the coordinates of each entry of T are equal. Then the coordinates of the first $r-1$ entries are equal. However, since $S_r' = S_{r-1}\mathsf{e}_r^{\mathsf{F}_{r-1}}$ and by the hypothesis that $\Sigma_{\le |S_{r-1}|-1}(S_{r-1}) \cap H_{r-1} = \emptyset$, this is only possible if either $T \mid \mathsf{e}_r^{\mathsf{F}_{r-1}}$ or $S_{r-1} \mid T$. In the former case, since $r \ge 2$, it is clear that $\sigma(T) \notin H$. In the latter case, since $\sigma(S_{r-1}) = \mathsf{F}_{r-1}\mathsf{e}_1 + \cdots + \mathsf{F}_{r-1}\mathsf{e}_{r-1}$ and $S_r' = S_{r-1}\mathsf{e}_r^{\mathsf{F}_{r-1}}, \sigma(T) \in H$ only if $T = S_r'$. Thus $\Sigma_{\le |S_r'|-1}(S_r') \cap H_r = \emptyset$ follows, completing the proof. $\qquad\square$

Theorem 3.13. *Let* $r \in \mathbb{N}_{\ge 2}$, *let* $(\mathsf{e}_1, \ldots, \mathsf{e}_r)$ *denote the standard basis of* $\mathbb{Z}^r$, *and let* $G_0 = G_r^+ \cup -G_r^+$. *Then*

$$\mathsf{F}_{r+2} \le \mathsf{D}(G_0) \le \frac{(2^r - r - 1)}{2^r}(r+2)^{\frac{r+2}{2}} \le (r+2)^{\frac{r+2}{2}}.$$

Proof. The upper bounds follow from Corollary 3.11 (see the discussion following the corollary). It remains to show $\mathsf{F}_{r+2} \le \mathsf{D}(G_0)$. Let $S_0 \in \mathcal{F}(G_r^+)$ be a sequence satisfying the conclusions of Proposition 3.12 and define

$$U = S_r(-\mathsf{e}_1 - \cdots - \mathsf{e}_r)^{\mathsf{F}_r} \in \mathcal{F}(G_0).$$

Then $|U| = |S_r| + \mathsf{F}_r = \mathsf{F}_{r+1} + \mathsf{F}_r = \mathsf{F}_{r+2}$ and $\sigma(U) = \sigma(S_r) - (\mathsf{F}_r\mathsf{e}_1 + \cdots + \mathsf{F}_r\mathsf{e}_r) = 0$. Moreover, by the definition of U it is clear that if $T \mid U$ is a zero-sum subsequence with $T = T^+ T^-$ where $\mathrm{supp}(T^+) \subset G_r^+$ and $\mathrm{supp}(T^-) \subset -G_r^+$, then we must have $\sigma(T^+) \in H$ with $T^+ \mid S_r$. However, since $\sigma(S_r) \in H_r$ and $\Sigma_{\le |S_r|-1}(S_r) \cap H_r = \emptyset$, this is only possible if either T^+ is trivial or $T^+ = S_r$. If T^+ is trivial, then $T = (-\mathsf{e}_1 - \cdots - \mathsf{e}_r)^{|T|}$ which is a zero-sum sequence only if T is trivial. If $T^+ = S_r$, then $\sigma(T^+) = \sigma(S_r) = \mathsf{F}_r\mathsf{e}_1 + \cdots + \mathsf{F}_r\mathsf{e}_r$ and it then follows from the definition of U that the only way T can be a zero-sum is if $T = U$. Therefore, U is a zero-sum sequence of length F_{r+2} having no proper nontrivial zero-sum subsequences. Thus $\mathsf{D}(G_0) \ge \mathsf{F}_{r+2}$ as desired. $\qquad\square$

Remark 3.14. Let $r \in \mathbb{N}, (\mathsf{e}_1, \ldots, \mathsf{e}_r)$, and G_0 be as in Theorem 3.13 but with $r > 1$. From Theorem 3.7 we know that $\mathsf{D}^{\mathrm{elm}}(G_0) \le \mathsf{D}(G_0) \le \frac{1}{2}|G_0|\mathsf{D}^{\mathrm{elm}}(G_0)$. But the expected value of $\mathsf{D}^{\mathrm{elm}}(G_0)$ is much larger than $|G_0|$, meaning that $\mathsf{D}(G_0) \approx \mathsf{D}^{\mathrm{elm}}(G_0)$. Indeed, by a computer based search we have verified that $\mathsf{D}(G_0) = \mathsf{D}^{\mathrm{elm}}(G_0)$ for $r \in \{2, 3\}$. However, whether $\mathsf{D}(G_0) = \mathsf{D}^{\mathrm{elm}}(G_0)$ remains true for $r \ge 4$ is not known. For the module-theoretic relevance of this question see [2, Lemma 6.9 and Corollary 6.10].

4. Monoids of Modules Over Commutative Noetherian Local Rings

In this section, we study direct-sum decompositions of finitely generated modules over commutative Noetherian local rings. With R a commutative Noetherian local ring and $\mathcal{C}$ a class of finitely generated modules over R, it is well-known that the monoid of modules $\mathcal{V}(\mathcal{C})$ is Krull. Our first lemma summarizes some basic information about $\mathcal{C}$.

Lemma 4.1. *Let $(R, \mathfrak{m})$ be a commutative Noetherian local ring with maximal ideal $\mathfrak{m}$ and let $(\widehat{R}, \widehat{\mathfrak{m}})$ denote its $\mathfrak{m}$-adic completion. Let $\mathcal{C}$ denote the class of all finitely generated R-modules and let $\mathcal{C}'$ denote a subclass of $\mathcal{C}$ such that $\mathcal{V}(\mathcal{C}') \subset \mathcal{V}(\mathcal{C})$ is a divisor-closed submonoid.*

(1) *The embedding $\mathcal{V}(\mathcal{C}) \hookrightarrow \mathcal{V}(\widehat{\mathcal{C}})$, defined by $[M] \mapsto [M \otimes \widehat{R}]$, is a divisor homomorphism into the free abelian monoid $\mathcal{V}(\widehat{\mathcal{C}})$, where $\widehat{\mathcal{C}}$ denotes the class of finitely generated $\widehat{R}$-modules. In particular, $\mathcal{V}(\mathcal{C})$ is a Krull monoid.*

(2) *$\mathcal{V}(\mathcal{C}')$ is a Krull monoid whose class group is an epimorphic image of a subgroup of the class group of $\mathcal{V}(\mathcal{C})$. If $\mathcal{V}(\mathcal{C}')$ is tame, then each of the arithmetical finiteness results of Proposition 2.1 hold.*

(3) *Suppose, in addition, that R is one-dimensional and reduced (no nonzero nilpotent elements). Let G denote the class group of $\mathcal{V}(\mathcal{C}')$ and let $G_P \subset G$ denote the set of classes containing prime divisors.*

   (a) *The class group G of $\mathcal{V}(\mathcal{C}')$ is a finitely generated abelian group.*

   (b) *$\mathcal{V}(\mathcal{C}')$ is tame if and only if $\mathsf{D}(G_P) < \infty$ if and only if G_P is finite.*

Proof. The monoid $\mathcal{V}(\widehat{\mathcal{C}})$ is free abelian by the theorem of Krull–Remak–Schmidt–Azumaya. Also, the embedding defined by $[M] \mapsto [M \otimes \hat{R}]$ is a divisor homomorphism by [53] or [40, Corollary 1.15], and hence $\mathcal{V}(\mathcal{C})$ is a Krull monoid, proving (1).

Let $H = \mathcal{V}(\mathcal{C})$ and suppose that the inclusion $H \hookrightarrow F = \mathcal{F}(P)$ a divisor theory. Then $\mathcal{C}(H) = \mathsf{q}(F)/\mathsf{q}(H)$. If $H' \subset H$ is divisor-closed, the inclusion $H' \hookrightarrow F' = \mathcal{F}(P')$, where $P' = \{p \in P : p \text{ divides some } a \in H' \text{ in } F\}$, is a divisor homomorphism. Note that $H' = F' \cap \mathsf{q}(H')$ and that $\mathsf{q}(H') = \mathsf{q}(F') \cap \mathsf{q}(H)$. Therefore, the homomorphism $\mathsf{q}(F') \to \mathsf{q}(F)/\mathsf{q}(H)$, defined by $a \mapsto a\mathsf{q}(H)$ for each $a \in \mathsf{q}(F')$, has kernel $\mathsf{q}(F') \cap \mathsf{q}(H) = \mathsf{q}(H')$. Consequently, there exists a monomorphism $\mathsf{q}(F')/\mathsf{q}(H') \to \mathsf{q}(F)/\mathsf{q}(H) = \mathcal{C}(H)$. Finally, [28, Theorem 2.4.7] implies that the class group of H' is an epimorphic image of a subgroup of $\mathsf{q}(F')/\mathsf{q}(H')$, and hence of a subgroup of $\mathcal{C}(H)$, proving (2).

We first consider the first statement of (3). The class group of $\mathcal{V}(\mathcal{C})$ is free abelian of finite rank by [38, Theorem 6.3]. Therefore, by statement (2), the class group G of $\mathcal{V}(\mathcal{C}')$ is finitely generated. We now consider the second statement of (3). Since G is finitely generated, [30, Theorem 4.2] implies that $\mathcal{V}(\mathcal{C}')$ is tame if and only if $\mathsf{D}(G_P) < \infty$. Again, since G is finitely generated, [28, Theorem 3.4.2] implies that $\mathsf{D}(G_P) < \infty$ if and only if G_P is finite. $\qquad\square$

Remark 4.2. We now make two brief remarks on certain hypothesis in statement (3) of Lemma 4.1.

(1) The assumption that R is reduced can be slightly weakened (see [38, Sec. 6]). However, the assumption that R is one-dimensional is essential for guaranteeing that the class group of $\mathcal{V}(\mathcal{C}')$ is finitely generated (see [40, Lemma 2.16]). Indeed, there is a two-dimensional complete Noetherian local Krull domain S whose class group is not finitely generated. By a result of Heitman, there is a factorial two-dimensional Noetherian local domain R whose completion is isomorphic to S. Then, if $\mathcal{C}$ is the class of finitely generated torsion-free R-modules, the class group of $\mathcal{V}(\mathcal{C})$ is isomorphic to the class group of S (see [2, Sec. 5]), and hence not finitely generated.

(2) Both characterizations in (3)(b) strongly depend on the fact that the class group of $\mathcal{V}(\mathcal{C}')$ is finitely generated and thus the hypothesis that R is one-dimensional in (3)(a) is critical for the results of this section pertaining to monoids of modules.

Let R be a one-dimensional reduced commutative Noetherian local ring and let $\mathcal{C}'$ be a class of finitely generated R-modules such that $\mathcal{V}(\mathcal{C}')$ is a divisor-closed submonoid of the monoid $\mathcal{V}(\mathcal{C})$ of all finitely generated R-modules. Then $\mathcal{V}(\mathcal{C}')$ is a Krull monoid with finitely generated class group G and we let $G_P \subset G$ denote the set of classes containing prime divisors. By Proposition 2.3, sets of lengths in $\mathcal{V}(\mathcal{C}')$ can be studied in the monoid $\mathcal{B}(G_P)$ of zero-sum sequences over G_P. Using the combinatorial results of Sec. 3, we study sets of lengths in such Krull monoids in Theorem 4.3 in the case where G_P contains the set $G_r^+ \cup G_r^-$ where G_r^+ denotes the nonzero vertices of the hypercube in a finitely generated free abelian group. We then investigate finer arithmetical invariants of $\mathcal{V}(\mathcal{C}')$ for small values of r. We conclude this section with an explicit example of a monoid of modules realizing such a Krull monoid (see Corollary 4.7).

Theorem 4.3. *Let H be a Krull monoid whose class group $\mathcal{C}(H)$ is free abelian with basis $(e_1, \ldots, e_r)$ for some $r \geq 2$, and let $G_P \subset \mathcal{C}(H)$ denote the set of classes containing prime divisors. Suppose that G_P is finite and that $G_P = -G_P$ with $G_P \supset G_r^+$.*

(1) *There exists $M \in \mathbb{N}_0$ such that, for each $k \geq 2$, $\mathcal{U}_k(H) = L' \cup L^* \cup L''$, where L^* is an interval, $L' \subset \min L^* + [-M, -1]$, and $L'' \subset \max L^* + [1, M]$.*
(2) *For each $k \in \mathbb{N}$,*
$$\rho_{2k}(H) = k\mathsf{D}(G_P) \geq k\mathsf{F}_{r+2} \quad and$$
$$k\mathsf{D}(G_P) + 1 \leq \rho_{2k+1}(H) \leq k\mathsf{D}(G_P) + \frac{\mathsf{D}(G_P)}{2}.$$

Upper bounds for $\mathsf{D}(G_P)$, and hence for each $\rho_k(H)$, then follow from Theorem 3.13.

Proof. By Proposition 2.3, $\mathcal{B}(G_P)$ and H are tame and thus, by Lemma 2.2, it suffices to prove all assertions about H for the monoid $\mathcal{B}(G_P)$.

By Proposition 2.1(3), it suffices to show that $\min \Delta(G_P) = 1$. If $U = e_1 e_2 (-e_1 - e_2)$, then $U \in \mathcal{A}(G_P)$ with $\mathsf{L}((-U)U) = \{2, 3\}$. Therefore $1 \in \Delta(G_P)$, proving (1).

Since $G_P \supset G_r^+ \cup -G_r^+$, the monoid $\mathcal{B}(G_r^+ \cup -G_r^+)$ is a divisor closed submonoid of $\mathcal{B}(G_P)$ and thus $\mathsf{D}(G_P) \geq \mathsf{D}(G_r^+ \cup -G_r^+)$. Theorem 3.13 now implies that $\mathsf{D}(G_P) \geq \mathsf{D}(G_r^+ \cup - G_r^+) \geq \mathsf{F}_{r+2}$. The inequalities involving $\rho_k(G_P)$ now follow easily from the definitions. Indeed, if $k \in \mathbb{N}$ and $A \in \mathcal{B}(G_P)$ with

$$A = 0^{\mathsf{v}_0(A)} U_1 \cdot \ldots \cdot U_k = 0^{\mathsf{v}_0(A)} V_1 \cdot \ldots \cdot V_l,$$

where $U_1, \ldots, U_k, V_1, \ldots, V_l \in \mathcal{A}(G_P) \backslash \{0\}$ are minimal zero-sum sequences, then

$$2l \leq \sum_{i=1}^{l} |V_i| = |A| - \mathsf{v}_0(A) = \sum_{i=1}^{k} |U_i| \leq k\mathsf{D}(G_P),$$

whence $l \leq k\mathsf{D}(G_P)/2$ and thus $\rho_k(G_P) \leq k\mathsf{D}(G_P)/2$. If $U = g_1 \cdot \ldots \cdot g_l$ is a minimal zero-sum sequence of length $|U| = l = \mathsf{D}(G_P)$, then

$$U^k (-U)^k = \prod_{i=1}^{l} ((-g_i) g_i)^k \tag{25}$$

and hence $\rho_{2k}(G_P) \geq k\mathsf{D}(G_P)$. Multiplying each side of (25) by any fixed atom, we find that $\rho_{2k+1}(G_P) \geq k\mathsf{D}(G_P) + 1$. $\qquad \square$

We now make a conjecture that claims that, at least for sufficiently large $k \in \mathbb{N}$, the first statement of Theorem 4.3 holds with $M = 0$.

Conjecture 4.4. *Let H be a Krull monoid as in Theorem 4.3 and suppose, in addition, that $G_P = G_r^+ \cup -G_r^+$. Then there exists $k^* \in \mathbb{N}$ such that for each $k \geq k^*$, $\mathcal{U}_k(H)$ is an interval.*

We now discuss one possible strategy for proving Conjecture 4.4. If one could show that there exists $A^* \in \mathcal{B}(G_P)$ with $\mathsf{L}(A^*)$ an interval and $\max \mathsf{L}(A^*)/\min \mathsf{L}(A^*) = \mathsf{D}(G_P)/2$, then there must exist $k^* \in \mathbb{N}$ such that for each $k \geq k^*$, $\mathcal{U}_k(H)$ is an interval (see [20, Theorem 3.1]). Unfortunately, this strategy seems to require knowledge of the precise value of the Davenport constant, which is currently known only for $r \in [2, 3]$. Even for $r = 4$, it seems to be computationally infeasible to compute the Davenport constant. However, for small values of r we can provide a direct proof of Conjecture 4.4. Indeed, we are even able to show in Corollary 4.6 that the conjecture holds when $r \in [2, 3]$ for $k^* = 2$. We first provide a simple lemma.

Lemma 4.5. *Let G be a free abelian group of rank $r \in \mathbb{N}$, let $(e_1, \ldots, e_r)$ denote a basis for G, and set $G_0 = G_r^+ \cup -G_r^+$. For every $U \in \mathcal{A}(G_0)$ with $|U| \geq 3$ and any $g \in G_0$, we have $\mathsf{v}_g(U) < |U|/2$.*

Proof. Let $U = g^k g_1 \cdot \ldots \cdot g_l \in \mathcal{A}(G_0)$ where $k = \mathsf{v}_g(U)$ and $|U| = k + l$. By symmetry we may suppose that $g \in G_r^+$. Clearly, we have $k \leq l = |U| - k$ and hence

$k \leq |U|/2$. Assume for the sake of contradiction that $2k = |U|$. Then $g_1, \ldots, g_l \in -G_r^+$. For each $h \in G_0$ with $h = \sum_{i=1}^r a_i e_i$ for some $a_1, \ldots, a_r \in \mathbb{Z}$, we denote by $S(h)$ the set $S(h) = \{i \in [1, r] : a_i \neq 0\}$. We now obtain that

$$S(g_1) \cup \cdots \cup S(g_l) \subset S(g) \subset \bigcap_{i=1}^l S(g_i),$$

and hence $g_1 = \cdots = g_l = -g$. Unless $U = g(-g)$, this contradicts the fact that $U \in \mathcal{A}(G_0)$. □

Corollary 4.6. *Let H be a Krull monoid whose class group G is free abelian with basis $(e_1, \ldots, e_r)$ for some $r \in [2, 3]$. Let $G_P \subset G$ denote the set of classes containing prime divisors and suppose that $G_P = G_r^+ \cup -G_r^+$.*

(1) *Suppose $r = 2$. Then $\mathsf{c}(H) = \mathsf{t}(H) = \mathsf{D}(G_P) = 3$ and, for each $k \geq 2$, the set $\mathcal{U}_k(H)$ is an interval. Moreover, for all $k \geq 2$ and $j \in [0, 1], \rho_{2k+j}(H) = 3k + j$.*
(2) *Suppose $r = 3$. Then $\mathsf{c}(H) = \mathsf{D}(G_P) = 5$ and, for each $k \geq 2$, the set $\mathcal{U}_k(H)$ is an interval. Moreover, for all $k \geq 2$ and $j \in [0, 1], \rho_{2k+j}(H) = 5k + j$.*

Proof. When $r = 2$, the statement $\mathsf{c}(H) = \mathsf{t}(H) = \mathsf{D}(G_P) = 3$ is a simple observation. Indeed, this setting is a special case of [2, Theorem 6.4]. The remaining assertions follow from Theorem 4.3.

We now assume $r = 3$. A lengthy technical proof or a computer search shows that $\mathsf{D}(G_P) = 5$ and that $\{V \in \mathcal{A}(G_P) : |V| = 5\} = \{V_i, -V_i : i \in [1, 4]\}$ where

$$V_1 = (e_1 + e_2)(e_1 + e_3)(e_2 + e_3)(-e_1 - e_2 - e_3)^2,$$

$$V_2 = (e_1 + e_2)(e_1 + e_3)(-e_2 - e_3)(-e_1)^2,$$

$$V_3 = (-e_1 - e_3)(e_1 + e_2)(e_2 + e_3)(-e_2)^2, \quad \text{and}$$

$$V_4 = (-e_1 - e_2)(e_1 + e_3)(e_2 + e_3)(-e_3)^2.$$

Lemma 4.5 then implies that each minimal zero-sum sequence of length four is squarefree. Clearly, $\mathsf{L}((-V_1)V_1) = \{2, 5\}$. Proposition 2.3 now implies that $5 \leq 2 + \max \Delta(H) \leq \mathsf{c}(H) \leq \mathsf{D}(G_P) = 5$, whence $\mathsf{c}(H) = 5$. Again, by Proposition 2.3, it suffices to prove the remaining assertions for $\mathcal{B}(G_P)$.

Theorem 4.3 implies that $\rho_{2k}(G_P) = k\mathsf{D}(G_P) = 5k$ and that $5k + 1 \leq \rho_{2k+1}(G_P) \leq 5k + 2$ for all $k \in \mathbb{N}$. In order to prove that $\mathcal{U}_k(H)$ is an interval for each $k \in \mathbb{N}$, it suffices to show that $\mathcal{U}_k(G_P) \cap \mathbb{N}_{\geq k}$ is an interval for each $k \in \mathbb{N}$. Indeed, this follows from a simple symmetry argument (see [20, Lemma 3.5]). We proceed to conclude the proof by proving the following three claims **A1**, **A2**, and **A3** which clearly imply the assertion.

A1. For each $k \in \mathbb{N}$, $\mathcal{U}_{2k}(G_P) \cap \mathbb{N}_{\geq 2k} = [2k, 5k]$.
A2. For each $k \in \mathbb{N}$, $\rho_{2k+1}(G_P) = 5k + 1$.
A3. For each $k \in \mathbb{N}$, $\mathcal{U}_{2k+1}(G_P) \cap \mathbb{N}_{\geq 2k+1} = [2k + 1, 5k + 1]$.

Proof of A1. For any $k \in \mathbb{N}$, the inclusion $\mathcal{U}_{2k}(G_P) \cap \mathbb{N}_{\geq 2k} \subset [2k, \rho_{2k}(G_P)] = [2k, 5k]$ is clear. To prove the reverse inclusion we proceed by induction on k. For every $j \in [3, 5]$ there exists $U \in \mathcal{A}(G_P)$ with $|U| = j$, and it follows that $\{2, j\} \subset \mathsf{L}((-U)U)$. Therefore, $\mathcal{U}_2(G_P) = [2, 5]$ and, together with the induction hypothesis, we see that for $k \geq 2$,

$$[2k, 5k] = [2, 5] + [2k - 2, 5k - 5] = \mathcal{U}_2(G_P) + \mathcal{U}_{2k-2}(G_P) \subset \mathcal{U}_{2k}(G_P).$$

Proof of A2. Assume for the sake of contradiction that there exists $k' \in \mathbb{N}$ such that $\rho_{2k'+1}(G_P) = 5k' + 2$, and let $k \in \mathbb{N}$ denote the smallest integer with this property. Let $B \in \mathcal{B}(G_P)$ be such that

$$B = A_1 \cdot \ldots \cdot A_{2k+1} = B_1 \cdot \ldots \cdot B_{5k+2}, \quad \text{where}$$

$$A_1, \ldots, A_{2k+1}, B_1, \ldots, B_{5k+2} \in \mathcal{A}(G_P).$$

Then

$$10k + 4 \leq \sum_{\nu=1}^{5k+2} |B_\nu| = |B| = \sum_{\nu=1}^{2k+1} |A_\nu| \leq 5(2k + 1),$$

and hence, after renumbering if necessary, $|A_1| = \cdots = |A_{2k}| = 5, A_{2k+1} \in [4, 5], |B_1| = \cdots = |B_{5k+1}| = 2$, and $|B_{5k+2}| \in [2, 3]$. Suppose there are $i, j \in [1, 2k + 1]$ with $i < j$ such that $A_i = -A_j$. Without loss of generality, $i = 1$ and $j = 2$ in which case $B' = A_3 \cdot \ldots \cdot A_{2k+1}$ satisfies $2k - 1, 5(k - 1) + 2 \in \mathsf{L}(B')$, contradicting the minimality of k.

Thus there are distinct $U_1, \ldots, U_4 \in \{V_i, -V_i : i \in [1, 4]\}$ such that $\{A_i : i \in [1, 2k]\} \subset \{U_1, \ldots, U_4\}$ with $U_i \neq -U_j$ for all $i, j \in [1, 4]$. Suppose that $U_1 = g_1^2 g_2 g_3 g_4$ such that $|\{i \in [1, 2k] : A_i = U_1\}| \geq \lceil k/2 \rceil$. Then $\mathsf{v}_{g_1}(B) \geq \max\{k, 2\}$. By inspection of all $W \in \mathcal{A}(G_P)$ with $|W| = 5$, it follows that $\mathsf{v}_{-g_1}(U_1 \cdot \ldots \cdot U_4) = 0 = \mathsf{v}_{-g_1}(A_1 \cdot \ldots \cdot A_{2k})$, and hence $\mathsf{v}_{-g_1}(B) = \mathsf{v}_{-g_1}(A_{2k+1})$. Since $|B_1| = \cdots = |B_{5k+1}| = 2, |B_{5k+2}| \in [2, 3]$, and $\mathsf{v}_{g_1}(B_{5k+2}) \leq 1$, it follows that

$$\mathsf{v}_{-g_1}(A_{2k+1}) = \mathsf{v}_{-g_1}(B) \geq \mathsf{v}_{g_1}(B) - 1 \geq \max\{k - 1, 1\},$$

and thus $|A_{2k+1}| = 4$. Therefore $|B_{5k+2}| = 2, \mathsf{v}_{-g_1}(A_{2k+1}) = \mathsf{v}_{-g_1}(B) = \mathsf{v}_{g_1}(B) \geq \max\{k, 2\}$. However, as we noted before, each minimal zero-sum sequences of length four is squarefree, a contradiction.

Proof of A3. Let $k \in \mathbb{N}$. By assertion **A2**, $\mathcal{U}_{2k+1}(G_P) \cap \mathbb{N}_{\geq 2k+1} \subset [2k + 1, \rho_{2k+1}(G_P)] = [2k + 1, 5k + 1]$. On the other hand,

$$[2k + 1, 5k + 1] = 1 + [2k, 5k] = \mathcal{U}_1(G_P) + \mathcal{U}_{2k}(G_P) \subset \mathcal{U}_{2k+1}(G_P). \qquad \square$$

Module theory provides an abundance of examples of Krull monoids satisfying the assumptions of Theorem 4.3 (see [2, Sec. 4], or [40, Chap. 1] for various realization results). Below we provide one specific example of a Krull monoid where the

set G_P of classes containing prime divisors is precisely $G_P = G_r^+ \cup G_r^-$ where G_r is the set of nonzero vertices of the r-dimensional hypercube.

Corollary 4.7. *Let $(R, \mathfrak{m})$ be a one-dimensional analytically unramified commutative Noetherian local domain with unique maximal ideal $\mathfrak{m}$. Further assume that the $\mathfrak{m}$-adic completion $\widehat{R}$ of R has $r + 1$ minimal primes $\mathfrak{q}_1, \ldots, \mathfrak{q}_{r+1}$, and let M be a torsion-free R-module whose completion $\widehat{M} = M \otimes_R \widehat{R}$ satisfies*

$$\widehat{M} \cong \bigoplus_{\emptyset \neq I \subset [1, r+1]} \frac{\widehat{R}}{\bigcap_{i \in I} \mathfrak{q}_i}.$$

Then $\mathcal{V}(\mathrm{add}(M))$ is a Krull monoid with class group $G \cong \mathbb{Z}^r$ and the set G_P of classes containing prime divisors satisfies $G_P = G_r^+ \cup -G_r^+$, where G, G_P, and G_r^+ are as in Theorem 4.3. Therefore, all arithmetical invariants, including all sets $\mathcal{U}_k(\mathcal{V}(\mathrm{add}(M)))$, are as described in Theorem 4.3 and Corollary 4.6.

Proof. The statements about G and G_P follow from [2, Example 4.21]. The arithmetical consequences then follow from Theorem 4.3 and Corollary 4.6. $\square$

5. Monoids of Modules Over Prüfer Rings

In this section, we study classes $\mathcal{C}$ of finitely presented modules over Prüfer rings and characterize the algebraic structure of the monoid $\mathcal{V}(\mathcal{C})$. Specifically, we study certain classes of projective modules over various types of Prüfer rings, and show that they are always half-factorial. We also study the catenary and tame degrees of these monoids. We first recall the definition of a Prüfer ring and related topics as well as that of finitely primary monoids. For a general reference on modules over Prüfer rings, the reader may wish to consult the monograph of Fuchs and Salce [22]. For modules over Prüfer rings with zero-divisors, we refer the reader to [18]. For additional information on finitely primary monoids, see [28, Secs. 2.7, 3.1, and 4.3].

A *Prüfer ring* is a commutative ring in which every finitely generated regular ideal is invertible. A commutative ring R has

- the $1\frac{1}{2}$ *generator property* if, for any invertible ideal $I \subset R$ and any regular element $a \in I \backslash \mathrm{rad}(R)I$, there exists an element $b \in R$ such that $I = Ra + Rb$,
- *small zero-divisors* if for every zero-divisor $a \in R$ and any ideal $A \subset R$, $A + aR = R$ implies that $A = R$.

A monoid H is called *finitely primary* if there exist $s, \alpha \in \mathbb{N}$ with the following properties:

H is a submonoid of a factorial monoid $F = F^\times \times [p_1, \ldots, p_s]$ for s pairwise non-associated prime elements $p_1, \ldots, p_s$ satisfying

$$H \backslash H^\times \subset p_1 \cdot \ldots \cdot p_s F \quad \text{and} \quad (p_1 \cdot \ldots \cdot p_s)^\alpha F \subset H.$$

In this case, we say that H is finitely primary of *rank s* and *exponent α*.

It is easy to show that the complete integral closure of such a finitely primary monoid H is F, and hence H is a Krull monoid if and only if H is factorial. Moreover, H_{red} is finitely generated if and only if $s = 1$ and $(\widehat{H}^{\times} : H^{\times}) < \infty$. The main (and motivating) examples of finitely primary monoids stem from ring theory. For example, if R is a one-dimensional local Mori domain with nonzero conductor $(R : \widehat{R})$ and $\widehat{R}$ denotes the complete integral closure of R, then $R^{\bullet}$ is finitely primary (see [28, Proposition 2.10.7]). The arithmetic of finitely primary monoids is well-studied (see [28, Secs. 2.7, 3.1, and 4.3]). In particular, the sets $\mathcal{U}_k(H)$ are finite (for one $k \geq 2$, or equivalently for all $k \geq 2$) if and only if $s = 1$. In our main results of this section, Theorems 5.1 and 5.3, we apply these arithmetical results to monoids of modules of modules over Prüfer rings.

Theorem 5.1. *Let R be a Prüfer ring such that R has the $1\frac{1}{2}$ generator property and R has small zero-divisors. Let $\mathcal{C}_{\mathrm{proj}}$ be the class of finitely generated projective R-modules.*

(1) *$\mathcal{V}(\mathcal{C}_{\mathrm{proj}})$ is a finitely primary monoid of rank 1 and of exponent 1. Moreover, $\mathcal{V}(\mathcal{C}_{\mathrm{proj}})$ is finitely generated if and only if $\mathrm{Pic}(R)$ is finite.*
(2) *If $\mathcal{C}_{\mathrm{proj}}$ does not satisfy KRSA, then $\mathcal{V}(\mathcal{C}_{\mathrm{proj}})$ is half-factorial with $\mathsf{c}(\mathcal{V}(\mathcal{C}_{\mathrm{proj}})) = \mathsf{t}(\mathcal{V}(\mathcal{C}_{\mathrm{proj}})) = 2$.*

Proof. We first consider statement (1). By [18, Corollary 4], every module P in $\mathcal{C}_{\mathrm{proj}}$ is isomorphic to $R^{n-1} \oplus I$ where $n \in \mathbb{N}$ is the rank of P, I is an invertible ideal, and the isomorphism class of P is determined by that of I and by the rank n. Thus the map $\varphi : \mathcal{V}(\mathcal{C}_{\mathrm{proj}}) \to H = (\mathrm{Pic}(R) \times \mathbb{N}) \cup \{(0,0)\}$, defined by $P \mapsto ([I], n)$, is an isomorphism. By definition, $H \subset \mathrm{Pic}(R) \times (\mathbb{N}_0, +)$ is finitely primary of rank 1 and exponent 1. Since H is reduced, we obtain that $(\widehat{H}^{\times} : H^{\times}) = |\mathrm{Pic}(R)|$ and thus H is finitely generated if and only if $\mathrm{Pic}(R)$ is finite.

We now compute the catenary and tame degrees based on the monoid described in (1). Suppose that $\mathcal{C}_{\mathrm{proj}}$ does not satisfy KRSA. Then H is not factorial and hence $\mathsf{c}(H) \geq 2$. By [28, Theorem 3.1.5], every finitely primary monoid of rank 1 and exponent α satisfies $\mathsf{c}(H) \leq \mathsf{t}(H) \leq 3\alpha - 1$, and hence $\mathsf{c}(H) = \mathsf{t}(H) = 2$. Now Proposition 2.1 implies that $\Delta(H) = \emptyset$, that is, H is half-factorial. $\square$

We now restrict our attention to direct-sum decompositions of modules over Prüfer domains. Specifically, we consider the class of all finitely presented modules, including torsion modules. Before proceeding, we recall the following characterization of Prüfer domains. An integral domain R is Prüfer if and only if the following two equivalent conditions are satisfied:

(P1) The torsion submodule of any finitely generated R-module M is isomorphic to a direct summand of M (see [22, Chap. V, Corollary 2.9]).

(P2) Every finitely generated R-module is projective if and only if it is torsion-free (see [22, Chap. V, Theorem 2.7]).

In Theorem 5.3 we can provide a more precise arithmetical description of $\mathcal{V}(\mathcal{C})$, where $\mathcal{C}$ is the class of finitely presented modules over a Prüfer domain, if we further assume that the domain is h-local. We now recall this class of integral domains. Let R be a domain and, for an ideal $I \subset R$, let $\Omega(I)$ denote the set of maximal ideals of R containing I. Note, that $|\Omega(I)| = 1$ implies that R/I is local and hence indecomposable as an R-module. Also recall that a domain R has finite character if each nonzero element of R is contained in at most finitely many maximal ideals of R. Now, we say that an integral domain R is h-*local* if the following equivalent conditions are satisfied (see [44, Theorem 2.1]):

(H1) R has finite character and each nonzero prime ideal of R is contained in a unique maximal ideal.
(H2) For each nonzero ideal I of R, R/I has a decomposition $\bigoplus_{\nu=1}^m R/I_\nu$ with $|\Omega(I_1)| = \cdots = |\Omega(I_m)| = 1$.
(H3) Each torsion R-module M is canonically isomorphic to $\bigoplus_{\mathfrak{p} \in \max(R)} M_\mathfrak{p}$.

For additional information on h-local domains we refer the reader to the survey paper by Olberding [44] and to the monograph by Fontana, Houston, and Lucas [19]. The next proposition gathers together the module-theoretic results necessary for the arithmetical results we present in Theorem 5.3. We would like to thank Bruce Olberding for the short proof of Proposition 5.2(1).

Proposition 5.2. *Let R be a commutative ring.*

(1) *Let $m \in \mathbb{N}$ and let $I, I_1, \ldots, I_m$ be ideals of R such that $R/I \cong R/I_1 \oplus \cdots \oplus R/I_m$. If I is finitely generated and projective as an R-module, then $I_1, \ldots, I_m$ are also finitely generated and projective as R-modules.*
(2) *Let $m, n \in \mathbb{N}$ and let $I_1, \ldots, I_m, J_1, \ldots, J_n$ be ideals of R. If $I_m \subset \cdots \subset I_1$, $J_n \subset \cdots \subset J_1$, and*

$$R/I_1 \oplus \cdots \oplus R/I_m \cong R/J_1 \oplus \cdots \oplus R/J_n,$$

then $m = n$ and $I_\nu = J_\nu$ for each $\nu \in [1, m]$.
(3) *Let R be an h-local Prüfer domain and let M be a finitely presented R-module. Then, as an R-module, M decomposes as*

$$M \cong R/I_1 \oplus \cdots \oplus R/I_m \oplus I_{m+1} \oplus \cdots \oplus I_n,$$

where $n \in \mathbb{N}_0, m \in [0, n], I_m \subset \cdots \subset I_1$ are proper invertible ideals of R, and $I_{m+1}, \ldots, I_n$ are invertible ideals of R.

Proof. Clearly, it is sufficient to prove that I_1 is finitely generated and projective. We set $J = I_2 \cap \cdots \cap I_m$ and, by [22, Lemma 1.1, Chap. V], obtain that $I = I_1 \cap J$

and $R = I_1 + J$. Therefore, there is a short exact sequence

$$0 \to I \to I_1 \oplus J \to R \to 0$$

of R-modules where the second map is the embedding and where the third map is given by $(x, y) \mapsto x - y$ for all $x \in I_1$ and all $y \in J$. This sequence splits, and thus $I_1 \oplus J \cong R \oplus I$. Now, since I is finitely generated and projective, so is I_1. This proves (1).

For the proofs of statements (2) and (3), see Proposition 2.10 and Theorem 4.12 in [22, Chap. V]. $\qquad\square$

We now state our main results about finitely presented modules over Prüfer domains.

Theorem 5.3. *Let R be a Prüfer domain, $\mathcal{C}$ the class of all finitely presented R-modules, $\mathcal{C}_{\mathrm{tor}}$ the class of finitely presented torsion modules, and $\mathcal{C}_{\mathrm{proj}}$ the class of finitely generated projective modules.*

(1) $\mathcal{V}(\mathcal{C}) = \mathcal{V}(\mathcal{C}_{\mathrm{proj}}) \times \mathcal{V}(\mathcal{C}_{\mathrm{tor}})$.
(2) *If R has the $1\frac{1}{2}$ generator property, then $\mathcal{V}(\mathcal{C}_{\mathrm{proj}})$ is finitely primary of rank 1 and exponent 1.*
(3) *Assume, in addition, that R is h-local. Then $\mathcal{V}(\mathcal{C}_{\mathrm{tor}})$ is free abelian and $\mathcal{V}(\mathcal{C})$ is half-factorial. Furthermore, $\mathcal{V}(\mathcal{C})$ is either factorial or $\mathsf{c}(\mathcal{V}(\mathcal{C})) = \mathsf{t}(\mathcal{V}(\mathcal{C})) = 2$.*

Proof. The proof of (1) follows immediately from **(P1)** and **(P2)**, and statement (2) follows immediately from Theorem 5.1.

We now prove statement (3), and we begin by showing that $\mathcal{V}(\mathcal{C}_{\mathrm{tor}})$ is free abelian. Since R has finite character, R has the $1\frac{1}{2}$ generator property. Let M be a finitely presented nonzero torsion R-module. We now argue that the module M has a decomposition as a direct sum of indecomposable finitely presented R-modules, and that such a decomposition is unique up to isomorphism. By Proposition 5.2(3), M has a decomposition

$$M \cong R/I_1 \oplus \cdots \oplus R/I_m,$$

with $m \in \mathbb{N}$ and invertible ideals $I_m \subset \cdots \subset I_1 \subsetneq R$. Property **(H2)** then implies that for each $\nu \in [1, m]$, R/I_ν has a direct-sum decomposition into indecomposables modules, each of the form R/J where $J \subset R$ and $|\Omega(J)| = 1$. Since each I_ν is invertible, each I_ν is finitely generated and projective and, by Proposition 5.2(1), the same is true for all ideals J with R/J occurring in the direct-sum decomposition of R/I_ν. Thus, after replacing the R/I_ν with direct sums of finitely generated indecomposable R-modules of the form R/J and then renaming, we may suppose that each R/I_ν is an indecomposable R-module, that each I_ν is an invertible ideal, and that $|\Omega(I_\nu)| = 1$ for each $\nu \in [1, m]$.

Let $M \cong C_1 \oplus \cdots \oplus C_n$ be any direct-sum decomposition of M into indecomposable finitely presented R-modules. Then Proposition 5.2 and **(H2)** imply that,

for each $\nu \in [1, n], C_\nu \cong R/J_\nu$ for some invertible ideal $J_\nu \subset R$ with $|\Omega(J_\nu)| = 1$. Therefore $M \cong R/J_1 \oplus \cdots \oplus R/J_n$. Let $\mathfrak{p}$ be a maximal ideal of R. Then

$$M_\mathfrak{p} \cong (R/I_1)_\mathfrak{p} \oplus \cdots \oplus (R/I_m)_\mathfrak{p} \cong (R/J_1)_\mathfrak{p} \oplus \cdots \oplus (R/J_n)_\mathfrak{p}$$

and, by **(H3)**, it suffices to prove uniqueness for the $R_\mathfrak{p}$-module $M_\mathfrak{p}$. Since the set of all ideals in the valuation domain $R_\mathfrak{p}$ form a chain, the uniqueness follows from Proposition 5.2(2).

Suppose now that $\mathcal{V}(\mathcal{C})$ is not factorial. By (1) and (2), $\mathcal{V}(\mathcal{C}) \cong F \times D$ where F is free abelian and where D is not factorial, but is finitely primary of rank 1 and exponent 1. Then

$$\mathsf{c}(F \times D) = \mathsf{t}(F \times D) = \mathsf{c}(D) = \mathsf{t}(D) = 2$$

and hence $F \times D$ is half-factorial. $\qquad\square$

Remark 5.4. Since a Noetherian Prüfer domain is precisely a Dedekind domain and since in the Noetherian setting the concepts of finitely presented and finitely generated modules coincide, Theorem 5.3 also describes the monoid of all finitely generated modules over a Dedekind domain. Of course, these results can be obtained even more simply from the classical results of Steinitz. In Sec. 6, we consider direct-sum decompositions of yet another class of rings that generalize Dedekind domains.

6. Monoids of Modules Over Hereditary Noetherian Prime Rings

In this final section, we study classes $\mathcal{C}$ of finitely generated right modules over HNP rings, a generalization of Dedekind prime rings (see [43, Sec. 5.7]). Module theory over HNP rings is carefully presented in the monograph of Levy and Robson [41], and it is on this work that this section is based. We begin with the arithmetical preparations necessary to state the main result in this section, Theorem 6.5. There we give a characterization of the monoid of stable isomorphism classes of finitely generated projective right modules over HNP rings and use this information to study its arithmetic.

Proposition 6.1. *Let H_0 and D be monoids and define*

$$H = H_0 \propto D = (H_0 \backslash H_0^\times) \times D \cup H_0^\times \times \{1_D\}.$$

Then H is a submonoid of $H_0 \times D$. If $D = \{1_D\}$ or H_0 is a group, then $H = H_0$. Suppose that $D \neq \{1_D\}$ and that H_0 is not a group.

(1) $H^\times = H_0^\times \times \{1_D\}$.
(2) $\mathsf{q}(H) = \mathsf{q}(H_0) \times \mathsf{q}(D), H \subset H_0 \times D$ *is not saturated, and $\widehat{H} = \widehat{H}_0 \times \widehat{D}$. In particular, H is not completely integrally closed and hence not a Krull monoid.*
(3) *The projection $\theta\colon H \to H_0, (a, d) \mapsto a$ is a transfer homomorphism.*
(4) *Suppose H_0 and D are atomic. Let $u \in \mathcal{A}(H_0)$ and let $d \in D$. Then*

$$\max\{\omega(H_0, u), \omega(D, d)\} \leq \omega(H, (u, d)) \leq \omega(H_0, u) + \omega(D, d) + \epsilon,$$

where $\epsilon = 1$ if u is prime and $d \in D^{\times}$, and $\epsilon = 0$ otherwise. In particular, $\omega(H, (u, d)) < \infty$ if and only if $\omega(H_0, u) < \infty$ and $\omega(D, d) < \infty$, and if D is not a group, then $\omega(H) = \infty$ and $\mathsf{t}(H) = \infty$.

(5) Suppose H_0 is atomic, D is a group, and let $d \in D$. If $u \in \mathcal{A}(H_0)$ is a prime element, then

$$\omega(H, (u, d)) = 2, \quad \tau(H, (u, d)) = 1 \quad and \quad \mathsf{t}(H, (u, d)H^{\times}) = 2.$$

If $u \in \mathcal{A}(H_0)$ is not a prime element, then

$$\omega(H, (u, d)) = \omega(H_0, u), \quad \tau(H, (u, d)) = \tau(H_0, u) \quad and$$

$$\mathsf{t}(H, (u, d)H^{\times}) = \mathsf{t}(H_0, uH_0^{\times}).$$

In particular,

$$\mathsf{t}(H) = \max\{2, \mathsf{t}(H_0)\}, \quad \omega(H) = \max\{2, \omega(H_0)\} \quad and$$

$$\tau(H) = \max\{1, \tau(H_0)\}.$$

(6) Suppose H_0 and D are atomic. Let $d \in D$ and $a \in H_0 \backslash (\mathcal{A}(H_0) \cup H_0^{\times})$. Then $\mathsf{c}_H((a, d)) = 0$ if

- $D^{\times} = \{1, d\}$, and $\mathsf{Z}_{H_0}(a) = \{(uH_0^{\times})^2\}$ for some $u \in \mathcal{A}(H_0)$, or
- D is reduced, $d = 1$ and $\mathsf{c}_{H_0}(a) = 0$, or
- D is reduced, $d \in \mathcal{A}(D)$ and $\mathsf{Z}_{H_0}(a) = \{(uH_0^{\times})^k\}$ for some $u \in \mathcal{A}(H_0)$ and $k \in \mathbb{N}_{\geq 2}$.

In any other case, $\mathsf{c}_H((a, d)) = \max\{2, \mathsf{c}_{H_0}(a)\}$. In particular,

$$\mathsf{c}(H) = \max\{2, \mathsf{c}(H_0)\}.$$

Before proceeding with the proof, the reader may find it useful to note that

$$H = H_0 \propto D = \{(h, d) \in H_0 \times D : h \in H_0^{\times} \text{ only if } d = 1_D\}.$$

Also, it will be convenient in the proof to introduce the following notation. For a monoid S and elements $a, b \in S$, we write $a \,\|\, b$ to denote that $a \,|\, b$ and $b \nmid a$, that is, a is a strict divisor of b.

Proof of Proposition 6.1. It is easily checked that H is a submonoid of $H_0 \times D$ and that $H = H_0$ if $H_0 = H_0^{\times}$ or $D = \{1_D\}$. Assume now that $D \neq \{1_D\}$ and $H_0 \neq H_0^{\times}$. Fix $a_0 \in H_0 \backslash H_0^{\times}$ and $d_0 \in D \backslash \{1_D\}$.

The proof of (1) is clear.

For statement (2), note that we have $\mathsf{q}(H) \subset \mathsf{q}(H_0) \times \mathsf{q}(D)$ and must show the reverse inclusion. Let $a, b \in H_0$ and $d, e \in D$. Then $aa_0, ba_0 \in H_0 \backslash H_0^{\times}$ and thus $(ab^{-1}, de^{-1}) = (aa_0, d)(ba_0, e)^{-1} \in \mathsf{q}(H)$. Therefore $\mathsf{q}(H_0) \times \mathsf{q}(D) \subset \mathsf{q}(H)$. To see that the inclusion $H \subset H_0 \times D$ is not saturated, note that $(1, d_0) \in \mathsf{q}(H) \cap (H_0 \times D)$, but $(1, d_0) \notin H$.

We now show that $\widehat{H} = \widehat{H}_0 \times \widehat{D}$. Since $H \subset H_0 \times D$, and since both monoids have the same quotient group, it is certainly true that $\widehat{H} \subset \widehat{H_0 \times D} = \widehat{H}_0 \times \widehat{D}$.

For the reverse inclusion, let $(x, y) \in \widehat{H}_0 \times \widehat{D}$. By definition, there exist $c \in H_0$ and $d \in D$ such that $(c, d)(x, y)^n \in H_0 \times D$ for all $n \in \mathbb{N}_0$. But then $ca_0 x^n \in H_0 \backslash H_0^\times$ and $dy^n \in D$ for all $n \in \mathbb{N}_0$ and thus $(ca_0, d)(x, y)^n \in H$. Therefore $(x, y) \in \widehat{H}$. Since $H \neq \widehat{H}$, H is not completely integrally closed and thus not a Krull monoid.

We now prove statement (3). Clearly θ is surjective and, by (1), $\theta^{-1}(H_0^\times) = H^\times$. Let $(a, d) \in H$ and $b, c \in H_0$ such that $\theta((a, d)) = bc$. To establish that θ is a transfer homomorphism, it will suffice to find $e, f \in D$ such that $(b, e), (c, f) \in H$ and $ef = d$ since, in this case, $(a, d) = (b, e)(c, f), \theta((b, e)) = b$ and $\theta((c, f)) = c$. If b is not a unit in H_0, then $e = d$ and $f = 1$ gives $(b, d), (c, 1) \in H$ with $ef = d$. Similarly, if c is not a unit, then $e = 1$ and $f = d$ gives $(b, 1), (c, d) \in H$ with $ef = d$. If both $b, c \in H_0^\times$, then $a \in H_0^\times$ and it necessarily follows that $d = 1$. Now $e = f = 1$ gives $(b, 1), (c, 1) \in H$.

To show (4) we first observe that, for $(a, d), (b, e) \in H, (a, d) \,|\, (b, e)$ if and only if $a \,|\, b$ and $d = e$ or $a \,\|\, b$ and $d \,|\, e$. We first prove the correctness of the lower bound and for this we may assume that $k = \omega(H, (u, d)) < \infty$. Suppose that $m, k \in \mathbb{N}$ with $k \leq m$ and that $u_1, \ldots, u_m \in \mathcal{A}(H)$ are such that $u \,|\, u_1 \cdot \ldots \cdot u_m$. Then $(u, d) \,|\, (u_1, d) \cdot \ldots \cdot (u_m, d)$, and hence there exists a subproduct of at most k elements that is divisible by (u, d), say $(u, d) \,|\, (u_1, d) \cdot \ldots \cdot (u_k, d)$. But then $u \,|\, u_1 \cdot \ldots \cdot u_k$ and consequently $\omega(H_0, u) \leq k$. Similarly, suppose $m, k \in \mathbb{N}$ with $k \leq m$ and suppose $d_1, \ldots, d_m \in \mathcal{A}(D)$ are such that $d \,|\, d_1 \cdot \ldots \cdot d_m$. Then $(u, d) \,|\, (u^2, d_1) \cdot \ldots \cdot (u^2, d_m)$, and so there is a subproduct of at most k elements that is divisible by (u, d), say $(u, d) \,|\, (u^2, d_1) \cdot \ldots \cdot (u^2, d_k)$. But then $d \,|\, d_1 \cdot \ldots \cdot d_k$ and so $\omega(D, d) \leq k$.

We now verify the upper bound, and for this we may assume that $k = \omega(H, u) < \infty$ and $l = \omega(D, d) < \infty$. Suppose that $m, k, l \in \mathbb{N}$ with $m \geq k + l + \epsilon$ and suppose $(u_1, d_1), \ldots, (u_m, d_m) \in \mathcal{A}(H)$ are such that $(u, d) \,|\, (u_1, d_1) \cdot \ldots \cdot (u_m, d_m)$. Then $u \,|\, u_1 \cdot \ldots \cdot u_m$ and $d \,|\, d_1 \cdot \ldots \cdot d_m$, whence there are subsets $I, J \subset [1, m]$ with $|I| \leq k$ and $|J| \leq l$ such that $u \,|\, \prod_{i \in I} u_i$ and $d \,|\, \prod_{j \in J} d_j$. Now $|I \cup J| \leq k + l$ and, after renumbering and possibly enlarging I and/or J, we have that $I \cup J = [1, k + l]$ with $u \,|\, u_1 \cdot \ldots \cdot u_{k+l}$ and $d \,|\, d_1 \cdot \ldots \cdot d_{k+l}$. If u is not a prime element, then $k \geq 2$. If d is not a unit then $l \geq 1$. In either of these cases, $k + l \geq 2$ and thus $u \,\|\, u_1 \cdot \ldots \cdot u_{k+l}$. Therefore, $(u, d) \,|\, (u_1, d_1) \cdot \ldots \cdot (u_{k+l}, d_{k+l})$, and consequently $\omega(H, (u, d)) \leq k + l$. On the other hand, if u is a prime element and d is a unit (i.e. $k + l = 1$), then $u \,\|\, u_1 \cdot \ldots \cdot u_{k+l+1}$ and hence $(u, d) \,|\, (u_1, d_1) \cdot \ldots \cdot (u_{k+l+1}, d_{k+l+1})$.

The claim that $\omega(H, (u, d)) < \infty$ if and only if $\omega(H_0, u) < \infty$ and $\omega(D, d) < \infty$ is clear from the now-verified inequalities. If D is not a group, then (since D is atomic) there exists some atom $d \in D$. Then $\omega(D, d^k) \geq k$ for all $k \in \mathbb{N}$, and thus $\omega(H, (u, d^k)) \geq k$ for all $u \in \mathcal{A}(H_0)$. This implies that $\omega(H) = \infty$ and therefore $\mathsf{t}(H) = \infty$ by Proposition 2.1.

We now prove statement (5) in which case we assume that D is a group. Since, for any atomic monoid S and any non-prime atom $u \in S, \mathsf{t}(S, uS^\times) = \max\{\omega(S, u), \tau(S, u) + 1\}$ (see Sec. 2), it will suffice to establish the claim for the ω- and τ-invariants.

First observe that since every element of D is a unit, if $(a, e), (b, f) \in H$ with $a \mid b$, then $(a, e) \mid (b, f)(c, g)$ for any $(c, g) \in H \backslash H^{\times}$, and if $a \parallel b$ then $(a, e) \mid (b, f)$. In particular, if $k \geq 2, u_1, \ldots, u_k \in \mathcal{A}(H_0)$, and $u \in \mathcal{A}(H_0)$, then $u \mid u_1 \cdot \ldots \cdot u_k$ if and only if $(u, e) \mid (u_1, e_1) \cdot \ldots \cdot (u_k, e_k)$ for any (equivalently, all) $e, e_1, \ldots, e_k \in D$.

Since $d \in D^{\times}$, we have that $\omega(D, d) = 0$ and thus statement (4) implies that $\omega(H, (u, d)) = \omega(H_0, u)$ if u is not a prime element, and $\omega(H, (u, d)) \in [1, 2]$ if u is a prime element (since then $\omega(H_0, u) = 1$). Moreover, $(u, d) \mid (u, dd_0)^2$ but $(u, d) \nmid (u, dd_0)$, and hence $\omega(H, (u, d)) \geq 2$ implying that if u is a prime element, then $\omega(H, (u, d)) = 2$.

A similar argument shows that $\tau(H, (u, d)) \geq 1$ for all $u \in \mathcal{A}(H_0)$. Suppose that u is a prime element and that $k \in \mathbb{N}_{\geq 2}, (u_1, d_1), \ldots, (u_k, d_k) \in \mathcal{A}(H)$, and $(u, d) \mid (u_1, d_1) \cdot \ldots \cdot (u_k, d_k)$. Then $u \mid u_1 \cdot \ldots \cdot u_k$ and thus u divides one of $u_1, \ldots, u_k$, say $u \mid u_1$. But then $(u, d) \mid (u_1, d_1)(u_2, d_2)$, showing $\tau(H, (u, d)) \leq 1$.

Recall the definition of the τ-invariant as the supremum of a certain set from Eq. (1) in Sec. 2. If $u \in \mathcal{A}(H_0)$ is not a prime element, then the supremum of this set is attained for $k \geq 2$, and since $\mathsf{L}_H((a, e)) = \mathsf{L}_{H_0}(a)$ for all $a \in H_0$ and $e \in D$ (by statement (3)), it is immediate that $\tau(H, (u, d)) = \tau(H_0, u)$. We note that if $k = 2$, then there may be factorizations in H that contribute elements to this set that are not already contributed by factorizations in H_0. However, if $k = 2$, then we necessarily have that $\min \mathsf{L}(u^{-1}a) = 1$ and thus the result is the same as for $k > 2$.

We assume now that H_0 and D are atomic and make use of the transfer homomorphism from statement (3) in order to prove statement (6). By Lemma 2.2(3)(b), we have that $\mathsf{c}_{H_0}(a) \leq \mathsf{c}_H((a, d)) \leq \max\{\mathsf{c}(a, \theta), \mathsf{c}_{H_0}(a)\}$. We first show that $\mathsf{c}(a, \theta) \leq 2$. For an element $(b, e) \in H$ we write $\overline{(b, e)} = (b, e)H^{\times}$ for its class in H_{red}. Suppose that $z = \overline{(u_1, d_1)} \cdot \ldots \cdot \overline{(u_k, d_k)}$ and $z' = \overline{(u_1, d_1')} \cdot \ldots \cdot \overline{(u_k, d_k')}$ with $k, k' \geq 2, u_1, \ldots, u_k \in \mathcal{A}(H_0)$, and $d_1, \ldots, d_k, d_1', \ldots, d_k' \in D$ are two factorizations of (a, d) lying in the same fiber of θ. Then $\overline{(u_1, d_1)} \cdot \ldots \cdot \overline{(u_{k-1}, d_{k-1}d_k)} \, \overline{(u_k, 1)}$ is also a factorization of (a, d) lying in the same fiber and

$$\mathsf{d}\big(\overline{(u_1, d_1)} \cdot \ldots \cdot \overline{(u_{k-1}, d_{k-1})} \, \overline{(u_k, d_k)}, \, \overline{(u_1, d_1)} \cdot \ldots \cdot \overline{(u_{k-1}, d_{k-1}d_k)} \, \overline{(u_k, 1)}\big) \leq 2.$$

Inductively, we find a 2-chain from z to $\overline{(u_1, d_1 \cdot \ldots \cdot d_k)} \, \overline{(u_2, 1)} \cdot \ldots \cdot \overline{(u_k, 1)}$. Similarly, we find a 2-chain from z' to $\overline{(u_1, d_1' \cdot \ldots \cdot d_k')} \, \overline{(u_2, 1)} \cdot \ldots \cdot \overline{(u_k, 1)}$ and hence a 2-chain from z to z'. This shows that $\mathsf{c}(a, \theta) \leq 2$.

Recall that $\mathsf{c}_H((a, d)) = 0$ is equivalent to (a, d) having a unique factorization in H and, in any other case, $\mathsf{c}_H((a, d)) \geq 2$. Therefore, to establish the remaining claims, it will suffice to show that in each of the cases listed, the element (a, d) has a unique factorization and that in any other case, (a, d) has at least two distinct factorizations. From the inequalities already proven, we then obtain that $\mathsf{c}_H((a, d)) = \max\{2, \mathsf{c}_{H_0}(a)\}$.

We first suppose that D is not reduced. Assume that there exists $\varepsilon \in D^{\times} \backslash \{1, d\}$ and let $k \in \mathbb{N}_{\geq 2}$ and $u_1, \ldots, u_k \in \mathcal{A}(H_0)$ with $a = u_1 \cdot \ldots \cdot u_k$. Then

$$\overline{(u_1, d)} \, \overline{(u_2, 1)} \cdot \ldots \cdot \overline{(u_k, 1)} \quad \text{and} \quad \overline{(u_1, d\varepsilon^{-1})} \, \overline{(u_2, \varepsilon)} \, \overline{(u_3, 1)} \cdot \ldots \cdot \overline{(u_k, 1)}$$

are two distinct factorizations of (a, d) since $(d\varepsilon^{-1}, \varepsilon) \neq \{(d, 1), (1, d)\}$. Therefore, in this case, $\mathsf{c}_H((a, d)) \geq 2$.

We now consider the case where $D^\times = \{1, d\}$ with $d \neq 1$. If $\mathsf{Z}_{H_0}(a) = \{(uH_0^\times)^2\}$ for some $u \in \mathcal{A}(H_0)$, then the unique factorization of (a, d) is $\overline{(u, d)}\,\overline{(u, 1)}$. The remaining cases are as follows.

- If a has two distinct factorizations, then so does (a, d).
- If a has a unique factorization represented by $u_1 \cdot \ldots \cdot u_k$ for some $k \in \mathbb{N}_{\geq 2}$ and $u_1, \ldots, u_k \in \mathcal{A}(H_0)$ with $u_1 \not\simeq u_2$, then

$$\overline{(u_1, d)}\,\overline{(u_2, 1)}\,\overline{(u_3, 1)} \cdot \ldots \cdot \overline{(u_k, 1)} \quad \text{and} \quad \overline{(u_1, 1)}\,\overline{(u_2, d)}\,\overline{(u_3, 1)} \cdot \ldots \cdot \overline{(u_k, 1)}$$

  are two distinct factorizations of (a, d).

- If $a \simeq u^k$ for some $k \in \mathbb{N}_{\geq 3}$ and $u \in \mathcal{A}(H_0)$, then $\overline{(u, d)}\,\overline{(u, 1)}^{k-1}$ and $\overline{(u, d)}^3\,\overline{(u, 1)}^{k-3}$ are two distinct factorizations of (a, d) (here we use $d^2 = 1$).

Now suppose that D is reduced. If $d = 1$ and, for some $k \in \mathbb{N}_{\geq 2}$ and $u_1, \ldots, u_k \in \mathcal{A}(H_0), a = u_1 \cdot \ldots \cdot u_k$ represents the unique factorization of a in H_0, then $\overline{(u_1, 1)} \cdot \ldots \cdot \overline{(u_k, 1)}$ is the unique factorization of (a, d) in H. If $d \in \mathcal{A}(D)$ and for some $u \in \mathcal{A}(H_0)$ and $k \in \mathbb{N}_{\geq 2}, a \simeq u^k$ represents the unique factorization of a in H_0, then $\overline{(u, d)}\,\overline{(u, 1)}^{k-1}$ is the unique factorization of (a, d) in H.

We now show that in each of the remaining cases there exist two distinct factorizations of (a, d) in H. This is clear if a has two distinct factorizations in H_0, and we may assume in the following that a has a unique factorization represented by $a = u_1 \cdot \ldots \cdot u_k$ for some $k \in \mathbb{N}_{\geq 2}$ and $u_1, \ldots, u_k \in \mathcal{A}(H_0)$.

- If $d \neq 1$ is not an atom, there must exist $d_1, d_2 \in D \setminus \{1\}$ such that $d = d_1 d_2$. Then $(a, d) = (u_1, d)(u_2, 1) \cdot \ldots \cdot (u_k, 1)$ and $(a, d) = (u_1, d_1)(u_2, d_2)(u_3, 1) \cdot \ldots \cdot (u_k, 1)$ give rise to two distinct factorizations of (a, d).
- If $d \in \mathcal{A}(D)$ and there exist two distinct factors in the factorization of a, say $u_1 \not\simeq u_2$, then $(a, d) = (u_1, d)(u_2, 1)(u_3, 1) \cdot \ldots \cdot (u_k, 1)$ and $(a, d) = (u_1, 1)(u_2, d)(u_3, 1) \cdot \ldots \cdot (u_k, 1)$ give rise to two distinct factorizations of (a, d).

Note that we have so far showed that $\mathsf{c}(H_0) \leq \mathsf{c}(H) \leq \max\{2, \mathsf{c}(H_0)\}$. We now show that $2 \leq \mathsf{c}(H)$. Since H_0 is atomic, but not a group, there exists $a \in H_0 \setminus H_0^\times$. Taking a power of a if necessary, we may further assume that $a \in H_0 \setminus \mathcal{A}(H_0)$. If D is not reduced, then $2 \leq \mathsf{c}((a, 1)) \leq \mathsf{c}(H)$. If D is reduced, then, since $D \neq \{1_D\}$ by assumption, there exists $d \in D \setminus D^\times$ and, again taking a power of a if necessary, we may assume $d \in D \setminus \mathcal{A}(D)$. Thus $2 \leq \mathsf{c}((a, d)) \leq \mathsf{c}(H)$. □

We note that if $H_0 = \mathbb{N}_0$ and D is a group, then $H = \mathbb{N}_0 \propto D$ as above is a finitely primary monoid of rank 1 and exponent 1 (as discussed at the beginning of Sec. 5). Also, if H_0, D and E are monoids, then $(H_0 \propto D) \propto E = H_0 \propto (D \times E)$. Before defining, in Definition 6.3, a monoid that will later be used to describe the monoid of stable isomorphism classes of finitely generated projective

right modules over an HNP ring, we give a simple lemma that will be used in the proof of Proposition 6.4.

Lemma 6.2. *Let G be an additive abelian group and let $G_0 \subset G$ be a subset.*

(1) *If $G_0 = \{g\}$ for some torsion element $g \in G$, then $\mathcal{A}(G_0) = \{g^{\mathrm{ord}(g)}\}$.*
(2) *Let $n \in \mathbb{N}_{\geq 2}$, let $(e_1, \ldots, e_{n-1})$ be a family of independent elements of G each of infinite order, and let*

$$G_0 = \{A_i e_i : i \in [1, n-1]\} \cup \left\{ -\sum_{i=1}^{n-1} B_i e_i \right\},$$

where $A_i, B_i \in \mathbb{N}$ with $\gcd(A_i, B_i) = 1$ for each $i \in [1, n-1]$. Then $\mathcal{A}(G_0) = \{U\}$ where

$$U = \prod_{i=1}^{n-1} (A_i e_i)^{\frac{B_i L}{A_i}} \left(-\sum_{i=1}^{n-1} B_i e_i \right)^L \quad \text{with } L = \mathrm{lcm}\{A_i : i \in [1, n-1]\}.$$

In particular, $\mathcal{B}(G_0)$ is factorial in each case.

Proof. If $\mathcal{A}(G_0) = \{S\}$, then S is a prime element, in which case $\mathcal{B}(G_0)$ is factorial.

Statement (1) is clear and we now prove statement (2). From the definition of L, we have $A_i \mid B_i L$. Moreover, $\sigma(U) = \sum_{i=1}^{n-1} (\frac{B_i L}{A_i} A_i - B_i L) e_i = 0 \in G$ which shows that $U \in \mathcal{B}(G_0)$. Suppose that $S = \prod_{i=1}^{n-1} (A_i e_i)^{k_i} (-\sum_{i=1}^{n-1} B_i e_i)^k \in \mathcal{B}(G_0)$ for some $k_1, \ldots, k_n, k \in \mathbb{N}_0$. Note that the $k_1, \ldots, k_n$ are uniquely determined by k and $\sigma(S) = 0$. Therefore, to establish that U is an atom of $\mathcal{B}(G_0)$, and in fact the unique atom, it will suffice to show that $L \mid k$. Since $\sigma(S) = 0$, we have that $A_i k_i = k B_i$ for each $i \in [1, n-1]$. Since $\gcd(A_i, B_i) = 1$, this implies $A_i \mid k$ and hence $L = \mathrm{lcm}\{A_i : i \in [1, n-1]\} \mid k$. □

Definition 6.3. Let Ω be a set containing a designated element 0 and let $\mathbf{c} \in \mathbb{Q}_{>0}^\Omega$ such that $c_0 = 1$ and $c_i \in \mathbb{N}$ for all but finitely many $i \in \Omega$. Define

$$\mathbb{N}_0^\Omega(\mathbf{c}) = \{\mathbf{x} \in \mathbb{N}_0^\Omega : x_0 > 0 \text{ and } |\mathrm{supp}_{\mathbb{Z}^\Omega}(\mathbf{x} - x_0 \mathbf{c})| < \infty\} \cup \{\mathbf{0}\},$$

that is, $\mathbb{N}_0^\Omega(\mathbf{c})$ consists of those vectors which are almost everywhere equal to a nonzero multiple (determined by the coordinate x_0) of $\mathbf{c}$, together with the vector $\mathbf{0}$. We write $\ell(\mathbf{x}) = x_0$. If $|\Omega| < \infty$, then $\mathbb{N}_0^\Omega(\mathbf{c}) \cong (\mathbb{N} \times \mathbb{N}_0^{\Omega \setminus \{0\}}) \cup \{\mathbf{0}\} \cong (\mathbb{N} \times \mathbb{N}_0^{(\Omega \setminus \{0\})}) \cup \{\mathbf{0}\}$. Let Λ be a (possibly empty) set of finite, pairwise disjoint subsets of $\Omega \setminus \{0\}$ each containing at least two elements and, for each $I \in \Lambda$, let $C_I = \sum_{i \in I} c_i$. We assume that $C_I \in \mathbb{N}$ for all $I \in \Lambda$. Define

$$\mathbb{N}_0^\Omega(\mathbf{c}, \Lambda) = \left\{ \mathbf{x} \in \mathbb{N}_0^\Omega(\mathbf{c}) : \sum_{i \in I} x_i = C_I \ell(\mathbf{x}) \text{ for all } I \in \Lambda \right\}.$$

By definition, $\mathbb{N}_0^\Omega(\mathbf{c})$ and $\mathbb{N}_0^\Omega(\mathbf{c}, \Lambda)$ are reduced submonoids of $(\mathbb{N}_0^\Omega, +)$. However, note that the inclusion $\mathbb{N}_0^\Omega(\mathbf{c}) \subset \mathbb{N}_0^\Omega$ is not saturated. Indeed, if $\mathbf{x}, \mathbf{y} \in \mathbb{N}_0^\Omega(\mathbf{c})$, then

x divides **y** if and only if $\mathbf{y} - \mathbf{x} \in \mathbb{N}_0^\Omega(\mathbf{c})$, that is, $\mathbf{x} \leq \mathbf{y}$ and either $x_0 < y_0$ or $\mathbf{x} = \mathbf{y}$.

Proposition 6.4. *Let $H = \mathbb{N}_0^\Omega(\mathbf{c}, \Lambda)$ with $\Omega, \mathbf{c}$ and Λ as in Definition 6.3. Then H is a saturated submonoid of $\mathbb{N}_0^\Omega(\mathbf{c})$ and the map $\ell : H \to (\mathbb{N}_0, +)$ is a transfer homomorphism. In particular, $\mathbf{x} \in H$ is an atom if and only if $\ell(\mathbf{x}) = 1$, and H is half-factorial.*

(1) *If $|\Omega| = 1$, then $H = \mathbb{N}_0$.*

(2) *Suppose $2 \leq |\Omega| < \infty$ and $\bigcup \Lambda = \Omega \backslash \{0\}$. Write $\Omega = [0, r]$ with $r \in \mathbb{N}$, and $\Lambda = \{I_1, \ldots, I_n\}$ with $n \in \mathbb{N}$ and $\Omega \backslash \{0\} = I_1 \uplus \cdots \uplus I_n$. Set $C_i = C_{I_i}$ for all $i \in [1, n]$.*

   (a) *The map $j : H \hookrightarrow \mathbb{N}_0^r, \mathbf{x} \mapsto (x_1, \ldots, x_r)$ is a divisor theory (note that x_0 is omitted), and hence H is a Krull monoid. Then $G = \mathbb{Z}^r / \mathsf{q}(j(H))$ is the divisor class group and $G_P = \{\mathbf{e}_i + \mathsf{q}(j(H)) : i \in [1, r]\} \subset G$ is the set of classes containing prime divisors. There is a monomorphism*

$$\varphi^* : G \to \begin{cases} \mathbb{Z}/C_1\mathbb{Z} & \text{if } n = 1, \\ \mathbb{Z}^{n-1} & \text{if } n \geq 2 \end{cases}$$

   with the following properties:

   If $n = 1$, then φ^* is an isomorphism, $\varphi^*(G_P) = \{1 + C_1\mathbb{Z}\}$, and the unique class in G_P contains precisely $|I_1|$ prime divisors. Let $n \geq 2$ and, for all $i \in [1, n-1]$, set $A_i = \frac{C_n}{\gcd(C_i, C_n)}$ and $B_i = \frac{C_i}{\gcd(C_i, C_n)}$. Then $\varphi^*(G)$ is a full rank subgroup of $\mathbb{Z}^{n-1}$ and

$$\varphi^*(G_P) = \{A_i \mathbf{e}_i : i \in [1, n-1]\} \cup \left\{ -\sum_{i=1}^{n-1} B_i \mathbf{e}_i \right\}.$$

   For each $i \in [1, n-1]$, the class mapped onto the element $A_i \mathbf{e}_i$ contains precisely $|I_i|$ prime divisors and the class mapped onto $-\sum_{i=1}^{n-1} B_i \mathbf{e}_i$ contains precisely $|I_n|$ prime divisors.

   (b) *An element $\mathbf{u} \in H$ is an atom if and only if $\sum_{i \in I} u_i = C_I$ for all (equivalently, any) $I \in \Lambda$.*

   (c) *We have that $\mathsf{c}(H) \leq 2$. Moreover, the monoid H is factorial if and only if $|\Lambda| = 1$ and $C_1 = 1$. In this case, the prime elements are precisely the $\mathbf{e}_0 + \mathbf{e}_i$ for $i \in [1, r]$. If H is not factorial, then H contains no prime elements.*

   (d) *For any atom $\mathbf{u} \in \mathcal{A}(H)$ we have*

$$\sum_{I \in \Lambda} \left(C_I - \min_{i \in I} u_i \right) \leq \omega(H, \mathbf{u}) \leq \sum_{I \in \Lambda} C_I.$$

   (e) *If H is not factorial, then $\mathsf{t}(H) = \omega(H) = \sum_{I \in \Lambda} C_I$.*

(3) *If $2 \leq |\Omega| < \infty$ but $\bigcup \Lambda \subsetneq \Omega \backslash \{0\}$, set $\Omega' = \bigcup \Lambda \cup \{0\}$. Let $\mathbf{c}' \in \mathbb{Q}_{>0}^{\Omega'}$ be defined by $c_i' = c_i$ for all $i \in \Omega'$ and let $H' = \mathbb{N}_0^{\Omega'}(\mathbf{c}', \Lambda)$. Then $H \cong H' \propto \mathbb{N}_0^{|\Omega| - |\Omega'|}$.*

For each $(\mathbf{u}, \mathbf{x}) \in \mathcal{A}(H)$,

$$\max\{\omega(H', \mathbf{u}), |\mathbf{x}|\} \le \omega(H, (\mathbf{u}, \mathbf{x})) \le \omega(H', \mathbf{u}) + |\mathbf{x}| + 1$$

and thus $\omega(H, (\mathbf{u}, \mathbf{x})) < \infty$ *and* $\mathsf{t}(H, (\mathbf{u}, \mathbf{x})) < \infty$. *Moreover,* H *is a half-factorial FF-monoid with* $\mathsf{c}(H) = 2$ *and* $\mathsf{t}(H) = \infty$. *Also,* H *is not a Krull monoid.*

(4) *If* $|\Omega| = \infty$, *then*

  (a) $|\mathsf{Z}_H(\mathbf{x})| = \infty$ *for all* $\mathbf{x} \in H \backslash (\mathcal{A}(H) \cup H^\times)$. *In particular,* H *is not a FF-monoid and is therefore not a submonoid of a free abelian monoid,*

  (b) $\mathsf{c}_H(\mathbf{x}) = 2$ *for all* $\mathbf{x} \in H \backslash (\mathcal{A}(H) \cup H^\times)$ *and*

  (c) $\omega(H, \mathbf{u}) = \tau(H, \mathbf{u}) = \mathsf{t}(H, \mathbf{u}) = \infty$ *for all* $\mathbf{u} \in \mathcal{A}(H)$.

Proof. Suppose that $\mathbf{x}, \mathbf{y} \in H$ such that $\mathbf{x}$ divides $\mathbf{y}$ in $\mathbb{N}_0^\Omega(\mathbf{c})$. That is, $\mathbf{y} - \mathbf{x} \in \mathbb{N}_0^\Omega(\mathbf{c})$. For each $I \in \Lambda$, we have $\sum_{i \in I} y_i - x_i = C_I(\ell(\mathbf{y}) - \ell(\mathbf{x})) = C_I \ell(\mathbf{y} - \mathbf{x})$, and thus $\mathbf{y} - \mathbf{x} \in H$. Therefore, the inclusion $H \subset \mathbb{N}_0^\Omega(\mathbf{c})$ is saturated. We now show that $\ell : H \to \mathbb{N}_0$ is a transfer homomorphism. Clearly this map is surjective and, since H is reduced, it will suffice to show: If $\mathbf{x} \in H$ and $\ell(\mathbf{x}) = k + l$ with $k, l \in \mathbb{N}_0$, then there exist $\mathbf{y}, \mathbf{z} \in H$ such that $\mathbf{x} = \mathbf{y} + \mathbf{z}$ and $\ell(\mathbf{y}) = k, \ell(\mathbf{z}) = l$.

If one of k or l is 0 we may, without restriction, assume that $l = 0$. Then $\mathbf{y} = \mathbf{x}$ and $\mathbf{z} = \mathbf{0}$ give the result. From now on, we assume that $k, l > 0$. Let $\Lambda' \subset \Lambda$ consist of those finitely many $I \in \Lambda$ with $I \cap \mathrm{supp}(\mathbf{x} - \ell(\mathbf{x})\mathbf{c}) \ne \emptyset$. For $i \in \bigcup \Lambda'$, let $y_i', z_i' \in \mathbb{N}_0$ be such that $y_i' + z_i' = x_i$ and such that, for all $I \in \Lambda'$, it holds that $\sum_{i \in I} y_i' = C_I k$ and $\sum_{i \in I} z_i' = C_I l$. Define $\mathbf{y}$ and $\mathbf{z} \in \mathbb{N}_0^\Omega$ by

$$y_i = \begin{cases} y_i' & \text{if } i \in \bigcup \Lambda', \\ kc_i & \text{if } i \in \Omega \backslash \bigcup \Lambda', \end{cases} \qquad z_i = \begin{cases} z_i' & \text{if } i \in \bigcup \Lambda', \\ lc_i & \text{if } i \in \Omega \backslash \bigcup \Lambda'. \end{cases}$$

Then $\mathbf{y}, \mathbf{z} \in \mathbb{N}_0^\Omega(\mathbf{c}, \Lambda)$ with $\ell(\mathbf{y}) = k, \ell(\mathbf{z}) = l$, and $\mathbf{x} = \mathbf{y} + \mathbf{z}$ as required.

Statement (1) is clear and thus we suppose that $2 \le |\Omega| < \infty$ and $\Omega \backslash \{0\} = \bigcup \Lambda$.

We first verify statement (2)(a). Since $\sum_{i \in I_1} x_i = C_1 x_0$ and $C_1 \ne 0, j$ is injective. We now check that j is a divisor homomorphism. Let $\mathbf{x}, \mathbf{y} \in H$ with $j(\mathbf{x}) \le j(\mathbf{y})$. If $\mathbf{x} = \mathbf{y}$, there is nothing to show. If $\mathbf{x} \ne \mathbf{y}$, then there necessarily exists $I \in \Lambda$ with $\sum_{i \in I} x_i < \sum_{i \in I} y_i$ and hence $x_0 < y_0$. Consequently, $\mathbf{x}$ divides $\mathbf{y}$ in H and thus j is a divisor homomorphism. To prove that j is a divisor theory, we need to show that each standard basis vector $\mathbf{e}_i \in \mathbb{N}_0^r$ is the greatest common divisor of a finite, nonempty set in the image of j. Let $i \in [1, r]$, let $I_0 \in \Lambda$ be such that $i \in I_0$, and let $i' \in I_0 \backslash \{i\}$. Define

$$\mathbf{x} = \mathbf{e}_0 + \mathbf{e}_i + (C_{I_0} - 1)\mathbf{e}_{i'} + \sum_{I \in \Lambda \backslash \{I_0\}} C_I \mathbf{e}_{\min I}, \quad \text{and}$$

$$\mathbf{y} = \mathbf{e}_0 + C_{I_0} \mathbf{e}_i + \sum_{I \in \Lambda \backslash \{I_0\}} C_I \mathbf{e}_{\max I}.$$

Then $\mathbf{x}, \mathbf{y} \in H$ and, recalling that $|I| \geq 2$ for all $I \in \Lambda$, $\gcd_{(\mathbb{N}_0^r, +)}(j(\mathbf{x}), j(\mathbf{y})) = \mathbf{e}_i$. Thus $j : H \hookrightarrow \mathbb{N}_0^r$ is a divisor theory, H is a Krull monoid, and the class group of H is $\mathsf{q}(\mathbb{N}_0^r)/\mathsf{q}(j(H)) = \mathbb{Z}^r/\mathsf{q}(j(H))$.

We now determine the structure of the divisor class group of H and determine the set of classes containing prime divisors. Suppose that $n = 1$ and define $\varphi^* : \mathbb{Z}^r \to \mathbb{Z}/C_1\mathbb{Z}$ by $\mathbf{x} \mapsto \sum_{i=1}^r x_i + C_1\mathbb{Z}$. We see immediately that φ^* is an epimorphism and $\ker(\varphi^*) = \mathsf{q}(j(H))$ follows easily. Therefore $G \cong \mathbb{Z}/C_1\mathbb{Z}$. Since $\varphi^*(\mathbf{e}_i) = 1 + C_1\mathbb{Z}$ for all $i \in [1, r]$, G_P is as claimed. Now suppose that $n \geq 2$ and let $\varphi^* : \mathbb{Z}^r \to \mathbb{Z}^{n-1}$ be defined by

$$\mathbf{x} \mapsto \left(A_1 \sum_{i \in I_1} x_i - B_1 \sum_{i \in I_n} x_i, A_2 \sum_{i \in I_2} x_i - B_2 \sum_{i \in I_n} x_i, \ldots, \right.$$
$$\left. A_{n-1} \sum_{i \in I_{n-1}} x_i - B_{n-1} \sum_{i \in I_n} x_i \right).$$

One easily checks that $\ker(\varphi^*) = \mathsf{q}(j(H))$, showing that G embeds into $\mathbb{Z}^{n-1}$ via φ^*. Considering the images of $\mathbf{e}_1, \ldots, \mathbf{e}_r$ under φ^*, the description of G_P given in statement (2)(a) follows, and we see that $\varphi^*(G)$ is a subgroup of full rank of $\mathbb{Z}^{n-1}$.

For statement (2)(b), note that an element $\mathbf{u} \in H$ is an atom if and only if $\ell(\mathbf{u}) = 1$, and this is the case if and only if $\sum_{i \in I} u_i = C_I$ for all (equivalently, any) $I \in \Lambda$.

Consider statement (2)(c). By Lemma 6.2, $\mathcal{B}(G_P)$ is factorial and hence $\mathsf{c}(G_P) = 0$. Thus Proposition 2.3 implies that $\mathsf{c}(H) \leq \max\{2, \mathsf{c}(G_P)\} \leq 2$. If $|\Lambda| = 1$ and $C_1 = 1$, then j is surjective. Thus $H \cong \mathbb{N}_0^r$, showing that H is factorial. The prime elements of $(\mathbb{N}_0^r, +)$ are simply the standard basis vectors $\mathbf{e}_i$ for $i \in [1, r]$, and their preimages under j are precisely the elements $\mathbf{e}_0 + \mathbf{e}_i \in H$ for $i \in [1, r]$. If $|\Lambda| > 1$ or $C_1 > 1$, then no atom is prime and thus H is not factorial. Let $\mathbf{u} \in \mathcal{A}(H)$. The lower bound given in (2)(d), which we will soon verify, implies that $\omega(H, \mathbf{u}) \geq 2$ unless $|\Lambda| = 1, C_1 = 2$, and $u_1 = u_2 = 1$. Note that $\min_{i \in I} u_i \leq \lfloor C_I/2 \rfloor$ for each $I \in \Lambda$ and thus, in this case, $\mathbf{u}$ is not a prime element. In the remaining case, one easily checks that again, $\mathbf{u}$ is not a prime element.

We now verify the bounds on the omega invariant as given in (2)(d). If $|\Lambda| = 1$ and $C_1 = 1$, then H is factorial and the inequalities hold trivially. We assume from now on that this is not the case and hence $\sum_{I \in \Lambda} C_I \geq 2$. We first show that $\omega(H, \mathbf{u}) \leq \sum_{I \in \Lambda} C_I$. Let $k \in \mathbb{N}_{\geq 2}$ and let $\mathbf{v}_1, \ldots, \mathbf{v}_k \in \mathcal{A}(H)$ be such that $\mathbf{u}$ divides $\sum_{i=1}^k \mathbf{v}_i$. If $J \subset [1, k]$, then $\mathbf{u}$ divides $\sum_{j \in J} \mathbf{v}_j$ if and only if $\mathbf{u} \leq \sum_{j \in J} \mathbf{v}_j$. Since $\mathbf{u} \leq \sum_{i=1}^k \mathbf{v}_i$ and $\sum_{i \in \Omega \setminus \{0\}} u_i = \sum_{I \in \Lambda} C_I$, we can recursively construct a subset $J \subset [1, k]$ of size at most $\sum_{I \in \Lambda} C_I$ such that $\mathbf{u} \leq \sum_{j \in J} \mathbf{v}_j$. This is done by adding, in each step, a vector $\mathbf{v}_j$ with $v_{j,i} > 0$ for some $i \in \Omega \setminus \{0\}$ for which $u_i < \sum_{j \in J} v_{j,i}$.

We now prove the lower bound. By renumbering the coordinates if necessary, we may assume $u_{\min I} = \min_{i \in I} u_i$ for all $I \in \Lambda$. For $i \in \Omega$, let $I \in \Lambda$ be the unique

set containing i and define

$$\mathbf{v}_i = \mathbf{e}_0 + \mathbf{e}_i + (C_I - 1)\mathbf{e}_{\min I} + \sum_{I' \in \Lambda \setminus \{I\}} C_{I'} \mathbf{e}_{\min I'} \in \mathcal{A}(H).$$

Since $\mathbf{u} = \sum_{i \in \Omega} u_i \mathbf{e}_i = \sum_{I \in \Lambda} \sum_{i \in I} u_i \mathbf{e}_i$, it is clear that $\mathbf{u}$ divides $\sum_{I \in \Lambda} \sum_{i \in I \setminus \{\min I\}} u_i \mathbf{v}_i$. For $I \in \Lambda$ with $C_I > 1$, this is clear from the definition of the $\mathbf{v}_i$'s. If $C_I = 1$, then $u_{\min I} = 0$ since $|I| \geq 2$, yet $\mathbf{u}$ does not divide any proper subsum. Therefore

$$\omega(H, \mathbf{u}) \geq \sum_{I \in \Lambda} \left[\left(\sum_{i \in I} u_i \right) - \min_{i \in I} u_i \right] = \sum_{I \in \Lambda} \left(C_I - \min_{i \in I} u_i \right).$$

We now prove statement (2)(e). By definition, $\omega(H) = \sup_{\mathbf{u} \in \mathcal{A}(H)} \omega(H, \mathbf{u})$. Since H is half-factorial, $\mathsf{t}(H) = \omega(H)$. The element

$$\mathbf{u} = \mathbf{e}_0 + \sum_{I \in \Lambda} C_I \mathbf{e}_{\min I} \in H$$

is an atom and the bounds in (2)(d) imply that $\omega(H, \mathbf{u}) = \sum_{I \in \Lambda} C_I$. Therefore, $\omega(H) \geq \sum_{I \in \Lambda} C_I$ and the upper bound again follows from (2)(d).

Consider statement (3) and let $D = \mathbb{N}_0^{|\Omega| - |\Omega_0|}$. The isomorphism $H \cong H' \propto D$ is immediate from the definitions. Because D is factorial, we have $\omega(D, \mathbf{x}) = |\mathbf{x}|$ for all $\mathbf{x} \in D$. Proposition 6.1(6) together with (2)(c) implies that $\mathsf{c}(H) = 2$, and Proposition 6.1(4) implies the remaining inequalities. As a submonoid of a free abelian monoid, H is an FF-monoid (see [28, Corollary 1.5.7]). However, since it is not completely integrally closed in its quotient group (by Proposition 6.1(2)), it cannot be a Krull monoid.

We suppose for the remainder of the proof that $|\Omega| = \infty$. For notational ease we assume $\mathbb{N}_0 \subset \Omega$.

Consider statement (4)(a). To show $|\mathsf{Z}_H(\mathbf{x})| = \infty$ it suffices to show that there are infinitely many atoms of H dividing $\mathbf{x}$. By definition of H, there exists a finite subset $\Omega' \subset \Omega$ having the property that $I \cap \Omega' \neq \emptyset$ already implies $I \subset \Omega'$ for all $I \in \Lambda$ and such that for each $i \in \Omega \setminus \Omega'$ it holds that $x_i = \ell(x) c_i$ and $c_i \in \mathbb{N}$. We may assume that $\mathbb{N}_0 \subset \Omega \setminus \Omega'$. Moreover, we may assume that for each $I \in \Lambda$ we have $|I \cap \mathbb{N}| \leq 1$ and, if $i \in I \cap \mathbb{N}$, then $c_i = \min_{i' \in I} c_{i'}$.

For each $i \in \mathbb{N}$, define $\mathbf{u}_i \in \mathbb{N}_0^{\Omega}$ as follows: If there exists $I \in \Lambda$ with $i \in I$, then let $i' \in I \setminus \{i\}$. Note also that $i' \in \Omega \setminus \Omega'$ by the choice of Ω'. We set $u_{i,i} = 0, u_{i,i'} = c_i + c_{i'}$, and $u_{i,j} = c_j$ for each $j \in \Omega \setminus (\Omega' \cup \{i, i'\})$. On the other hand, if i is not contained in any $I \in \Lambda$, we set $u_{i,i} = 0$ and $u_{i,j} = c_j$ for each $j \in \Omega \setminus (\Omega' \cup \{i\})$. In either case, we can choose $u_{i,j} \leq x_{i,j}$ for each $j \in \Omega'$ such that $\sum_{j \in I} u_{i,j} = C_I$ for all $I \in \Lambda$ with $I \subset \Omega'$. Therefore, $\sum_{j \in I} u_{i,j} = C_I$ for all $I \in \Lambda$ and $u_{i,j} = c_j$ for all but finitely many $j \in \Omega$. Therefore $\mathbf{u}_i \in \mathcal{A}(H)$. By construction, $\mathbf{u}_i \leq \mathbf{x}$. In the first case, due to the minimal choice of c_i, we have $u_{i,i'} = c_i + c_{i'} \leq 2c_{i'} \leq \ell(\mathbf{x}) c_{i'} = x_{i'}$, with the last equality holding since $i' \in \Omega \setminus \Omega'$. Since $\ell(\mathbf{u}_i) = 1 < \ell(\mathbf{x})$, $\mathbf{u}_i$ divides $\mathbf{x}$. Clearly these atoms are pairwise distinct. By what we have just shown, H is not

an FF-monoid and therefore cannot a submonoid of a free abelian monoid by [28, Corollary 1.5.7].

We now compute the catenary degree of H as stated in (4)(b). From (4)(a), every nonzero, non-atom element has at least two distinct factorizations and hence $\mathsf{c}_H(\mathbf{x}) \geq 2$ for all $\mathbf{x} \in H$. We now show that $\mathsf{c}_H(\mathbf{x}) \leq 2$ by projecting to the finite case. Suppose that $\mathbf{x} = \mathbf{u}_1 + \cdots + \mathbf{u}_k = \mathbf{v}_1 + \cdots + \mathbf{v}_k$ for some $k \in \mathbb{N}_{\geq 2}$ and that $\mathbf{u}_1, \ldots, \mathbf{u}_k, \mathbf{v}_1, \ldots, \mathbf{v}_k \in \mathcal{A}(H)$. Let $\Omega' \subset \Omega$ be the smallest subset containing

$$\{j \in \Omega : u_{i,j} \neq c_j \text{ or } v_{i,j} \neq c_j \text{ for some } i \in [1,k]\} \cup \{0\}$$

and having the property that whenever $I \cap \Omega' \neq \emptyset$ for $I \in \Lambda$, then already $I \subset \Omega'$. Observe that Ω' is finite and that $\mathbf{u}_{i,j} = \mathbf{v}_{i',j} = c_j$ for all $i, i' \in [1,k]$ and $j \in \Omega \backslash \Omega'$. Set $\Lambda' = \{I \in \Lambda : I \subset \Omega'\}$ and let $\mathbf{c}' \in \mathbb{N}_0^{\Omega'}$ be defined by $c_i' = c_i$ for all $i \in \Omega'$. Define $H' = \mathbb{N}_0^{\Omega'}(\mathbf{c}', \Lambda')$ and note that there is a canonical projection $\pi : H \to H'$. By the finiteness of Ω' we immediately have that $\mathsf{c}(H') \leq 2$. Thus there exists a sequence of factorizations $z_1', \ldots, z_l' \in \mathsf{Z}(H')$ of $\pi(\mathbf{x})$ with $z_1' = \pi(\mathbf{u}_1) \cdot \ldots \cdot \pi(\mathbf{u}_k), z_l' = \pi(\mathbf{v}_1) \cdot \ldots \cdot \pi(\mathbf{v}_k)$ and such that $\mathsf{d}(z_i', z_{i+1}') \leq 2$ for all $i \in [1, l-1]$. For any factor $\mathbf{w}'$ occurring in some z_j', we lift it to $\mathbf{w} \in H$ by setting $w_j = u_{1,j}$ for all $j \in \Omega \backslash \Omega'$. Then $\mathbf{w} \in \mathcal{A}(H)$ and we can lift the factorizations $z_1', \ldots, z_l'$ of $\pi(\mathbf{x})$ to factorizations $z_1, \ldots, z_l$ of $\mathbf{x}$. These factorizations form a sequence connecting the two factorizations of $\mathbf{x}$ that we began with and have the property that $\mathsf{d}(z_i, z_{i+1}) \leq 2$ for all $i \in [1, l-1]$.

Finally, we consider the infinitude of the invariants given in (4)(c). Let $\mathbf{u} \in \mathcal{A}(H)$ and $k \in \mathbb{N}_{\geq 2}$. We will show that $\omega(H, \mathbf{u}) \geq k$, as the other invariants are then equal to $\omega(H, u)$ by half-factoriality. Note that $\mathbf{u}$ is necessarily nonzero in all but finitely many coordinates and so we may assume $u_i \neq 0$ for $i \in [1, k]$. Moreover, since Ω is infinite and each $I \in \Lambda$ is finite, we may suppose that no two $i, j \in [1, k]$ belong to the same set $I \in \Lambda$. (For the following argument it would suffice that $I \not\subset [1, k]$.) Modifying finitely many coordinates of $\mathbf{u}$ as needed, we can therefore construct $\mathbf{u}_1, \ldots, \mathbf{u}_k \in \mathcal{A}(H)$ such that $u_{i,j} = u_j \delta_{i,j}$ for all $i, j \in [1, k]$, and $u_{i,j} \geq u_j$ for all $i \in [1, k]$ and $j \in \Omega \backslash [1, k]$. Then $\mathbf{u}$ divides $\mathbf{u}_1 + \cdots + \mathbf{u}_k$, but for all $j \in [1, k], u_j > u_{1,j} + \cdots + u_{k,j} - u_{j,j} = 0$, showing that $\mathbf{u}$ divides no proper subsum. $\qquad\square$

Now with the appropriate arithmetical results, we apply Propositions 6.1 and 6.4 to study direct-sum decompositions of modules over HNP rings. For a HNP ring R, we consider the class $\mathcal{C}$ of finitely generated right R-modules. In the noncommutative setting there are two invariants (the *genus* and the *Steinitz class*) which describe the stable isomorphism class of a finitely generated projective right R-module P. In general, however, the isomorphism class of P is determined by its stable isomorphism class only if $\mathrm{udim}(P) \geq 2$ and such a module is indecomposable if and only if $\mathrm{udim}(P) = 1$. Thus the forthcoming description of the direct-sum decomposition of finitely generated projective right R-modules is one where the indecomposable factors are determined up to stable isomorphism. If R has the additional property that any two finitely generated projective right R-modules that are

stably isomorphic are already isomorphic, then this result is a description up to isomorphism.

Let V and W be two simple right R-modules. Then W is called a *successor* of V and V is called a *predecessor* of W if $\mathrm{Ext}_R^1(V, W) \neq \mathbf{0}$. Let $\mathcal{W}$ be a set of representatives (of isomorphism classes) of the simple unfaithful right R-modules. We note that every $V \in \mathcal{W}$ is contained in a unique *tower* of R (see [41, Sec. 19]). A *tower* $\mathcal{T}$ is a finite set of simple right R-modules, ordered with respect to the successor relationship, and having the following structure: Every tower $\mathcal{T}$ is either cyclically ordered and each simple module in $\mathcal{T}$ is unfaithful, in which case we say that $\mathcal{T}$ is a *cycle tower*, or $\mathcal{T}$ is linearly ordered and only the first module (the only module in the tower not having a predecessor) is faithful, in which case $\mathcal{T}$ is said to be a *faithful tower*. The length of a tower is the number of distinct modules contained in it, and a tower is *nontrivial* if it contains more than one module.

We briefly recall the notions of rank, genus and the Steinitz class of a finitely generated projective right R-module, as these invariants are used to describe stable isomorphism classes of finitely generated projective right R-modules. We refer the reader to [41, Secs. 33 and 35] for additional details. Let P be a finitely generated projective right R-module and let $V \in \mathcal{W}$. Then $M = \mathrm{ann}(V)$ is a maximal ideal of R, R/M is simple artinian, and the *rank of M at V*, denoted by $\rho(P, V) \in \mathbb{N}_0$, is defined to be the length of the R/M-module P/PM. If $\mathcal{T}$ is a cycle tower, we set $\rho(P, \mathcal{T}) = \sum_{V \in \mathcal{T}} \rho(P, V)$. Let $\mathrm{modspec}(R)$ denote the set of isomorphism classes of all unfaithful simple right R-modules together with the trivial module $\mathbf{0}$. For a finitely generated projective right R-module P, we set $\Psi_V(P) = \rho(P, V)$ if V is an unfaithful simple right R-module, and $\Psi_{\mathbf{0}}(P) = \mathrm{udim}(P)$. Then

$$\Psi(P) = (\Psi_V(P))_{V \in \mathrm{modspec}(R)} \in \mathbb{N}_0^{\mathrm{modspec}(R)}$$

is called the *genus* of P. Two finitely generated projective right R-modules P and Q are *stably isomorphic* if there exists a finitely generated projective right R-module X such that $P \oplus X \cong Q \oplus X$. We denote by $[P]$ the stable isomorphism class of P. The direct-sum operation on modules induces the structure of a commutative semigroup on the set of stable isomorphism classes, and by $\mathbf{K}_0(R)$ we denote its quotient group. The genus Ψ induces a homomorphism $\Psi^+ : \mathbf{K}_0(R) \to \mathbb{Z}^{\mathrm{modspec}(R)}$, and $G(R) = \mathrm{Ker}(\Psi^+)$ is called the *ideal class group of R*. By choosing a base point set $\mathcal{B}$ of nonzero finitely generated projective right R-modules consisting of exactly one module in each genus, and such that $\mathcal{B}$ is closed, up to isomorphism, under direct sums, we can associate to any nonzero finitely generated projective right R-module P a class $\mathcal{S}(P) = [P] - [B] \in G(R)$, where $B \in \mathcal{B}$. We call $\mathcal{S}(P)$ the *Steinitz class* of P and we set $\mathcal{S}(\mathbf{0}) = \mathbf{0}$.

We are now ready to characterize the semigroup of stable isomorphism classes of finitely generated projective right modules over an HNP ring. We note that the characterization of factoriality in Theorem 6.5 was already obtained by Levy and Robson in [41, Theorem 39.5].

Theorem 6.5. *Let R be an HNP ring and let H be the semigroup of stable isomorphism classes of finitely generated projective right R-modules with operation induced by the direct sum of modules.*

(1) *Let Ω denote the set of isomorphism classes of all unfaithful simple right R-modules which are contained in a nontrivial tower, together with the trivial module, denoted by $\mathbf{0}$. For $V \in \Omega \backslash \{\mathbf{0}\}$, let $c_V = \frac{\rho(R,V)}{\mathrm{udim}(R)} \in \mathbb{Q}_{>0}$ and $c_{\mathbf{0}} = 1$. Finally, let Λ be the set of all nontrivial cycle towers. Then*

$$H \cong \mathbb{N}_0^{\Omega}(\mathbf{c}, \Lambda) \propto G(R)$$

with the product $\propto$ as defined in Proposition 6.1. In particular, H is half-factorial, $\mathsf{c}(H) \leq 2$ and the following are equivalent:

   (a) *H is factorial (i.e. R satisfies stable uniqueness).*
   (b) *$G(R) = \mathbf{0}$ and either R has no nontrivial towers (i.e. R is a Dedekind prime ring), or R has a unique nontrivial tower $\mathcal{T}$ which is a cycle tower and $\rho(R, \mathcal{T}) = \mathrm{udim}(R)$.*

(2) *If R has infinitely many nontrivial towers, then $\mathsf{t}(H, [U]) = \omega(H, [U]) = \infty$ for each indecomposable finitely generated projective right R-module U.*

(3) *Suppose that R has only finitely many nontrivial towers. Then $\mathsf{t}(H, [U])$, $\omega(H, [U]) < \infty$ for each indecomposable finitely generated projective right R-module U. If R has at least one nontrivial faithful tower, then $\mathsf{t}(H) = \omega(H) = \infty$. If R has no nontrivial faithful tower, only finitely many nontrivial cycle towers $\mathcal{T}_1, \ldots, \mathcal{T}_n$ with $n \in \mathbb{N}_0$, and H is not factorial, then*

$$\mathsf{t}(H) = \omega(H) = \frac{1}{\mathrm{udim}(R)} \sum_{i=1}^{n} \rho(R, \mathcal{T}_i).$$

Proof. It will suffice to establish that $H \cong \mathbb{N}_0^{\Omega}(\mathbf{c}, \Lambda) \propto G(R)$ as the remaining claims will then follow from Proposition 6.4. The genus and the Steinitz class are both additive on direct sums and thus give rise to a monoid homomorphism $H \to \mathbb{N}_0^{\mathrm{modspec}(R)} \propto G(R)$, $[P] \mapsto (\Psi(P), \mathcal{S}(P))$. We note that $\mathcal{S}(\mathbf{0}) = \mathbf{0}$ and that $\mathbf{0}$ is the only module of uniform dimension zero. The genus satisfies a number of necessary conditions, namely that $\Psi_V(P) = c_V \Psi_0(P)$ for almost all $V \in \mathrm{modspec}(R)$ and that it has standard rank at every cycle tower, that is, $\rho(P, \mathcal{T}) = \sum_{V \in \mathcal{T}} c_V \Psi_0(P)$ for all cycle towers $\mathcal{T}$. In particular, if V is contained in a trivial tower, then this is necessary a cycle tower and hence implies $\Psi_V(P) = c_V \Psi_0(P)$. Thus instead of $\Psi(P)$ we may consider $\Psi'(P) \subset \mathbb{N}_0^{\Omega}$ were we omit the components corresponding to unfaithful simple right R-modules that are contained in a trivial tower. We obtain a homomorphism

$$\Phi : H \to \mathbb{N}_0^{\Omega}(\mathbf{c}, \Lambda) \propto G(R), \quad [P] \mapsto (\Psi'(P), \mathcal{S}(P)).$$

The main theorem of Levy and Robson (see [41, Theorem 35.13]) implies that the genus and the Steinitz classes are independent invariants, and that up to the

stated conditions on the rank, all values can be obtained. In other words, Φ is an isomorphism. $\square$

Remark 6.6. (1) If R is such that each two stably isomorphic right R-modules are isomorphic, then $H = \mathcal{V}(\mathcal{C}_{\text{proj}})$. If this is not the case, then H still provides information on direct-sum decompositions of finitely generated projective modules with the summands determined up to stable isomorphism. Specifically, we have the following: Clearly, if $P = U_1 \oplus \cdots \oplus U_k$ for some $k \in \mathbb{N}$ and indecomposable right R-modules $U_1, \ldots, U_k$, then $[P] = [U_1] + \cdots + [U_k]$ is a factorization of $[P]$ in H. Let P be a right R-module with $\text{udim}(P) \geq 2$ (that is, P is neither the zero module nor indecomposable). If $[P] = [U_1] + \cdots + [U_k]$ for some $k \in \mathbb{N}_{\geq 2}$ and atoms $[U_1], \ldots, [U_k]$ of H, then $[P] = [U_1 \oplus \cdots \oplus U_k]$ and, since $\text{udim}(P) \geq 2$, [41, Theorem 34.6] implies that $P \cong U_1 \oplus \cdots \oplus U_k$. Therefore, a factorization of $[P]$ in H gives rise to one of P into indecomposables, with the stable isomorphism classes of the indecomposable summands determined by the factorization of $[P]$ in H.

(2) An HNP ring R has finitely many nontrivial towers if and only if it is a multichain idealizer from a Dedekind prime ring S (see [41, Proposition 30.5]). A sufficient (but not necessary) condition for there to be no nontrivial faithful towers is for R to be right (equivalently, left) bounded (see [41, Lemma 18.2]).

(3) Let R be an HNP ring and let $\mathcal{C}$ denote the class of finitely generated right R-modules. Then $\mathcal{V}(\mathcal{C}) \cong \mathcal{V}(\mathcal{C}_{\text{tor}}) \times \mathcal{V}(\mathcal{C}_{\text{proj}})$ and $\mathcal{V}(\mathcal{C}_{\text{tor}})$ is factorial.

Proof. By [41, Corollary 12.16(ii), Theorem 12.18], every finitely generated right R-module decomposes uniquely as a direct sum of a torsion module and a torsion-free module. The first has finite length and hence has a unique decomposition into indecomposables, while the second is projective. $\square$

A detailed description of $\mathcal{V}(\mathcal{C}_{\text{tor}})$ is given in [41, Sec. 41].

Acknowledgments

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