Bounding linear forms on a lattice away from zero

Lenny Fukshansky
Claremont McKenna College

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The basic problem

Let

\[ L_1(x_1, \ldots, x_n), \ldots, L_t(x_1, \ldots, x_n) \]

be linear forms in \( n \) variables, and let \( \Lambda \subset \mathbb{R}^n \) be a lattice.
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*Obtain a bound of the form*

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\min \{|L_1(x)|, \ldots, |L_t(x)| : x \in \Lambda \setminus \{0\}\} \geq C
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Of course, this is not always possible, but this general problem has many interesting applications. The purpose of this talk is to describe two different directions where some variations of this question naturally come up.
A consequence of Liouville’s inequality

Let $K$ be a number field of degree $d$ with real algebraic numbers

$$1 = \alpha_1, \alpha_2, \ldots, \alpha_d \in K$$

forming a $\mathbb{Q}$-basis for $K$. Let $L(x_1, \ldots, x_n) = \sum_{i=1}^{n} \alpha_i x_i$. 
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forming a $\mathbb{Q}$-basis for $K$. Let $L(x_1, \ldots, x_n) = \sum_{i=1}^{n} \alpha_i x_i$. Then for any $0 \neq x \in \mathbb{Z}^d$,

$$\|L(x)\| > \frac{1}{C|x|^d},$$

where $\| \|$ stands for distance to the nearest integer, $C$ is an explicit constant depending on $d$ and $\alpha_1, \ldots, \alpha_d$, and $\| \|$ is the sup-norm.
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In particular, this is an instance of Problem 1 for one linear form with lower bound depending on $x$. 
Kronecker’s Approximation Theorem

Let us write \{ \} for the fractional part function, and let 1, \theta_1, \ldots, \theta_t be \mathbb{Q}\text{-linearly independent real numbers for } t \geq 1.

Theorem 1 (L. Kronecker, 1884)

The sequence of points \( (\{ q\theta_1 \}, \ldots, \{ q\theta_t \}) \) as q ranges over \( \mathbb{Z} \) is dense in the torus \( \mathbb{R}^t / \mathbb{Z}^t \).
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There are several different results in the literature addressing this question. We present ours.
Effective version of Kronecker’s Theorem

Theorem 2 (F., Moshchevitin, 2017)

Let $1, \theta_1, \ldots, \theta_t \in \mathbb{Q} \cap \mathbb{R}$ be \( \mathbb{Q} \)-linearly independent with height

$$h(\theta_j) \leq H \ \forall \ 1 \leq j \leq t \text{ and } e = [\mathbb{Q}(\theta_1, \ldots, \theta_t) : \mathbb{Q}].$$

Let $\varepsilon > 0$. Then $\forall (a_1, \ldots, a_t) \in \mathbb{R}^t$ $\exists q \in \mathbb{Z} \setminus \{0\}$ with

$$|q| \leq \left(2^{et(e-1)}(t+1)^{3e-1}(t!)^{2e}H^{2e(t+1)-e}\right)\varepsilon^{-e+1},$$

such that

$$\|q\theta_j - a_j\| \leq \varepsilon, \ 1 \leq j \leq t.$$  \hspace{1cm} (2)
Proof of Theorem 2

Liouville inequality implies that for any \( \mathbf{0} \neq \mathbf{m} \in \mathbb{Z}^t \) with sup-norm \( |\mathbf{m}| \leq Y \),

\[
\|m_1\theta_1 + \cdots + m_t\theta_t\| \geq C_1 Y^{-e+1}
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Now a homogeneous-inhomogeneous transference principle implies that for every $(a_1, \ldots, a_t) \in \mathbb{R}^t$ there exists $q \in \mathbb{Z} \setminus \{0\}$ with $|q|$ appropriately bounded from above in terms of $C_1$ and $Y$ so that

$$\max_{1 \leq j \leq t} \|q \theta_j - a_j\| \leq C_2 Y^{-1},$$

where $C_2$ is again explicitly computable in terms of $C_1$. 
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where \( C_2 \) is again explicitly computable in terms of \( C_1 \).

Hence for a given \( \varepsilon > 0 \) we simply need to compute the appropriate \( Y \) so that \( C_2 Y^{-1} \leq \varepsilon \).
Density of images of linear forms

A quick corollary of Kronecker’s theorem is density of the image of $\mathbb{Z}^n$ in the torus $\mathbb{R}^t / \mathbb{Z}^t$ under a collection of linear forms

$$L_1(x_1, \ldots, x_n), \ldots, L_t(x_1, \ldots, x_n) \in \mathbb{R}[x_1, \ldots, x_n].$$

Specifically, if 1 and coefficients of $L_1, \ldots, L_t$ are $\mathbb{Q}$-linearly independent, then $\forall \varepsilon > 0, (a_1, \ldots, a_t) \in \mathbb{R}^t \exists x \in \mathbb{Z}^n$ such that

$$\|L_i(x) - a_i\| < \varepsilon \forall 1 \leq i \leq t.$$ 

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We effectively extend this result in two different ways:

1. allow for the approximating vector $x$ above to come from an algebraic lattice $\Lambda$, not just $\mathbb{Z}^n$, 

2. exclude vectors from a prescribed union $\mathbb{Z}$ of projective varieties or sublattices not containing this lattice, that is we are interested in approximation vectors $x \in \Lambda \setminus \mathbb{Z}$. This provides information on distribution of such approximating vectors.
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The setup

1. **The lattice:** \( \Lambda_K(\mathcal{M}) \subset \mathbb{R}^{wd} \) be the lattice of rank \( sd \), which is the image of an \( \mathcal{O}_K \)-module \( \mathcal{M} \subset K^w \) of rank \( s \) under the standard Minkowski embedding, where \( K \) is a number field of degree \( d \).
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3. The linear forms: Let 
   - $K_1 = K(\Lambda_K(\mathcal{M}))$, 
   - $1, b_{11}, \ldots, b_{t(wd)} \in \overline{\mathbb{Q}}$ linearly independent over $K_1$; let $h(B) = \max_{i,j} h(b_{ij})$, $\ell = [K(b_{11}, \ldots, b_{t(wd)}): \mathbb{Q}]$, 
   - $L_i(x_1, \ldots, x_{wd})$ linear forms with coefficients $b_{i1}, \ldots, b_{i(wd)}$ for all $1 \leq i \leq t$, 

Define also $K = \ell^2 (t+1) - \ell$ and $h(M) = \min\{h(\alpha^{(K+1)}_{sd-1}) \alpha \in K, \alpha_{M} \subset \mathcal{O}_w K\}$. 


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Define also \( \mathfrak{K} = \ell^2(t + 1) - \ell \) and

\[
h(M) = \min \left\{ h(\alpha)^{(\mathfrak{K}+1)sd-1} h(\alpha^{-1})^\mathfrak{K} : \alpha \in K, \alpha M \subseteq \mathcal{O}_K^w \right\}.
\]
Avoidance results - I

With this notation, we have our first effective version of Kronecker's theorem with avoidance conditions.

Theorem 3 (F., Moshchevitin, 2017)

Let \( \mathbf{a} = (a_1, \ldots, a_t) \in \mathbb{R}^t \) and \( \varepsilon > 0 \). There exist \( \mathbf{x} \in \Lambda_K(M) \setminus \mathcal{Z} \) such that

\[
\| L_i(\mathbf{x}) - a_i \| < \varepsilon
\]

and

\[
|\mathbf{x}| \leq C_3 \left( \det \Lambda_K(M) \right)^{\frac{t}{s} + 1} h(B)^{\frac{t}{s}} h(M) \varepsilon^{-\ell + 1},
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where \( C_3 \) is an explicit constant depending on \( K, t, \ell, s, M \).

One special case of this theorem is when \( \mathcal{Z} \) is a union of linear spaces, which means that the point \( x \) in question is in \( \Lambda_K(M) \) but outside of a union of sublattices of smaller rank than \( \Lambda_K(M) \). What if the rank of such sublattices is equal to the rank of \( \Lambda_K(M) \)? The next theorem addresses this situation.
Avoidance results - II

Theorem 4 (F., Moshchevitin, 2017)

Let \( \mathbf{a} = (a_1, \ldots, a_t) \in \mathbb{R}^t \) and \( \varepsilon > 0 \). Let \( m > 0 \) and

\[
\Gamma_1, \ldots, \Gamma_m \subset \Lambda_K(\mathcal{M})
\]

be proper sublattices of full rank and respective determinants \( D_1, \ldots, D_m \), and let \( D = D_1 \cdots D_m \). Then there exists \( \mathbf{x} \in \Lambda_K(\mathcal{M}) \setminus \bigcup_{i=1}^m \Gamma_i \) such that

\[
\|L_i(\mathbf{x}) - a_i\| < \varepsilon
\]

and

\[
|\mathbf{x}| \leq C_4 \left( \sum_{i=1}^m \frac{D}{D_i} \right)^{\frac{K}{K+2}} \left( \det \Lambda_K(\mathcal{M}) \right)^{-m+1} h(B)^{\frac{K}{K+2}} h(\mathcal{M}) \varepsilon^{-\ell+1},
\]

where \( C_4 \) is an explicit constant depending on \( K, t, \ell, s, m \).
Proof of Theorems 3 and 4

• We first construct a point \( y \in \Lambda_K(\mathcal{M}) \) of controlled sup-norm, which is outside of \( \mathcal{Z} \) or \( \bigcup_{i=1}^{m} \Gamma_i \), respectively.

1. In the first case, we use the classical Minkowski's Successive Minima Theorem and a version of Alon's Combinatorial Nullstellensatz.
2. In the second case we employ a recent result of Henk and Thiel on points of small norm in a lattice outside of a union of full-rank sublattices.

• We use \( y \) to construct an infinite sequence of points \( n_y \) satisfying the above conditions, and use our Theorem 2 to obtain a value of the index \( n \) (depending on \( \epsilon > 0 \)) for which the required inequalities on values of linear forms are satisfied.
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- We use $y$ to construct an infinite sequence of points $ny$ satisfying the above conditions, and use our Theorem 2 to obtain a value of the index $n$ (depending on $\varepsilon > 0$) for which the required inequalities on values of linear forms are satisfied.
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It should be remarked that the most important feature of approximation results such as our Theorems 3 and 4 is the exponent on $\varepsilon$ in the bounds for $|x|$. As we show, this exponent is the same as in the corresponding bound of the effective version of Kronecker’s theorem that we use.
A problem from compressed sensing

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The following optimization problem is motivated by applications in compressed sensing, an area of applied mathematics concerned with robust reconstruction of sparse signals from few measurements. Let

\[ \mathbb{Z}_m^d := \left\{ \mathbf{x} \in \mathbb{Z}^d : |\{ i : x_i \neq 0 \}| \leq m \right\}. \]

**Problem 2**

*Construct an* \( m \times d \) *matrix* \( A \) *with* \( m < d \) *such that*

\[ |A| := \max\{|a_{ij}| : 1 \leq i \leq m, 1 \leq j \leq d\} \leq C_1, \]

*and for every nonzero* \( \mathbf{x} \in \mathbb{Z}_m^d, \)

\[ \| A\mathbf{x} \| \geq C_2, \]

*where* \( C_1, C_2 > 0, \) *and* \( \| \| \) *now stands for Euclidean norm.*
Liouville’s inequality again

This is an optimization problem with the goal of making $d$ as large as possible in comparison to $m$. Using Liouville’s inequality as before, we can take $m = 1$, arbitrary $d$, and $A = (\alpha_1, \ldots, \alpha_d)$, with real algebraic $1 = \alpha_1, \ldots, \alpha_d$ forming a $Q$-basis for a number field of degree $d$. Then for any $0 \neq x \in \mathbb{Z}^d$, $|Ax| > C|x|^d$ for an explicit constant $C$ depending on $d$ and $\alpha_1, \ldots, \alpha_d$: here $x$ does not even have to be sparse. However, this lower bound depends on $x$ and $d$, which is not sufficiently good for applications. We want an absolute lower bound $C_2$ in Problem 2.
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Existence of such matrices

Theorem 5 (F., Needell, Sudakov – 2017)

There exist $m \times d$ integer matrices $A$ with $m < d$ and bounded $|A|$ such that for any nonzero $x \in \mathbb{Z}^d_m$,

$$\|Ax\| \geq 1.$$ \hspace{1cm} (3)

In fact, for sufficiently large $m$, there exist such matrices with

$$|A| = 1 \text{ and } d = 1.2938m,$$

and there also exist such matrices with

$$|A| = k \text{ and } d = \Omega(\sqrt{k}m).$$
Proof of Theorem 5

We use two powerful results of Bourgain, Vu and Wood (2010):

Let $M_m$ be an $m \times m$ random matrix whose entries are 0 with probability $1/2$ and $\pm 1$ with probability $1/4$ each. Then the probability that matrix $M_m$ is singular is at most $(1/2 - o(1))^m$. 
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We combine these with the union bound to obtain an estimate on $d$ in terms of $m$ that guarantees positive probability of all $m \times m$ submatrices of $A$ being non-singular.
How much bigger than $m$ can $d$ be?

Theorem 6 (F., Needell, Sudakov – 2017)

For any integers $m \geq 3$, $k \geq 1$ and $m \times d$ integer matrix $A$ with $|A| = k$ satisfying (3) for all $s \leq m$, we must have

$$d \leq (2k^2 + 2)(m - 1) + 1.$$
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**Theorem 7 (Konyagin – 2018)**

If \( m \geq 2(\log k + 1) \), then

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d < 144k(\log k + 1)m.
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Theorem 7 (Konyagin – 2018)

If $m \geq 2(\log k + 1)$, then

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Theorem 8 (Sudakov – 2018)

For sufficiently large $m$,

$$d = O(k\sqrt{\log k} m).$$
Question 2

What is the optimal upper bound on \( d \) in terms of \( k \) and \( m \)? In particular, is it true that \( d = O(km) \)?
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Here is a low-dimensional example.

Example 1

Let $m = 3$, $d = 6$, $k = 1$, and define a $3 \times 6$ matrix

$$A = \begin{pmatrix} 1 & 1 & 1 & 1 & 1 & 1 \\ -1 & -1 & -1 & -1 & 0 & 0 \\ 1 & 1 & 0 & 0 & -1 & -1 \end{pmatrix}.$$  

This matrix has $|A| = 1$ and any three of its columns are linearly independent. Then for $s \leq 3$ and any $x \in \mathbb{Z}_s^6$, $\|Ax\| \geq 1$.  

Algebraic matrices

While our Theorem 5 gives the image vector $Ax$ bounded away from 0, many of its coordinates may still be 0. A variation of our original Problem 1 would ask for a matrix $A$ with $Ax$ bounded away from 0 and nonzero coordinates.
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**Corollary 9 (F., Needell, Sudakov – 2017)**

Let $B$ be the $d \times m$-transpose of a matrix satisfying (3) as guaranteed by Theorem 5. Let $\theta$ be an algebraic integer of degree $m$, and let $\theta = \theta_1, \theta_2, \ldots, \theta_m$ be its algebraic conjugates. For each $1 \leq i \leq m$, let $\theta_i = (1 \theta_i \ldots \theta_i^{m-1})^\top$, compute the $d \times m$ matrix

$$B \left( \theta_1 \quad \ldots \quad \theta_m \right),$$

and let $A$ be its transpose. Then $|A| = O(|B|m)$, for any $x \in \mathbb{Z}_s^d$, $0 < s \leq m$, $\|Ax\| \geq \sqrt{m}$ and the vector $Ax$ has all nonzero coordinates.
Proof of Corollary 9

We use the AM-GM inequality:

$$\frac{1}{m} \|Ax\|^2 = \frac{1}{m} \sum_{i=1}^{m} |(B\theta_i)x|^2 \geq \left( \prod_{i=1}^{m} |(B\theta_i)x|^2 \right)^{1/m} = |N_K((B\theta_1)x)|^{2/m},$$

where $N_K$ is the field norm, which is $\geq$ since $(B\theta_1)x$ is an algebraic integer.
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\]

where \(\mathbb{N}_K\) is the field norm, which is \(\geq\) since \((B\theta_1)x\) is an algebraic integer.

Since \((B\theta_1)x \neq 0\), its algebraic conjugates, which are the rest of the coordinates of the vector \(Ax\) must all be nonzero.
**Algebraic matrix example**

Let $m = 3$, $d = 6$, and take $K = \mathbb{Q}(\theta)$, where $\theta = 2^{1/3}$, then

$$\theta = \begin{pmatrix} 1 \\ \theta \\ \theta^2 \end{pmatrix}.$$
Algebraic matrix example

Let \( m = 3, \ d = 6, \) and take \( K = \mathbb{Q}(\theta), \) where \( \theta = 2^{1/3}, \) then

\[
\theta = \begin{pmatrix} 1 \\
\theta \\
\theta^2 \end{pmatrix}.
\]

Let \( k = 1 \) and take \( B \) to be the transpose of the matrix from Example 1, i.e.

\[
B = \begin{pmatrix} 1 & 1 & 1 \\
1 & 1 & 0 \\
1 & 0 & 1 \\
1 & 0 & -1 \\
1 & -1 & 0 \\
1 & -1 & -1 \end{pmatrix}.
\]
Algebraic matrix example

The number field $K$ has three embeddings into $\mathbb{C}$, given by

$$
\theta \mapsto \theta, \quad \theta \mapsto \xi \theta, \quad \theta \mapsto \xi^2 \theta,
$$

where $\xi = e^{\frac{2\pi i}{3}}$ is a third root of unity, i.e. $\theta$ is mapped to roots of its minimal polynomial by injective field homomorphisms that fix $\mathbb{Q}$. 

$$
|A| \leq 3^{\frac{3}{2}} \quad \text{and} \quad \|Ax\| \geq \sqrt{3} \quad \text{for every} \quad x \in \mathbb{Z}^3.
$$
Algebraic matrix example

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$$

where $\xi = e^{\frac{2\pi i}{3}}$ is a third root of unity, i.e. $\theta$ is mapped to roots of its minimal polynomial by injective field homomorphisms that fix $\mathbb{Q}$.

Hence we get the following $3 \times 6$ matrix:

$$
A = \begin{pmatrix}
1+\theta+\theta^2 & 1+\theta & 1+\theta^2 & 1-\theta^2 & 1-\theta & 1-\theta-\theta^2 \\
1+\xi \theta+\xi^2 \theta^2 & 1+\xi \theta & 1+\xi^2 \theta^2 & 1-\xi^2 \theta^2 & 1-\xi \theta & 1-\xi \theta-\xi^2 \theta^2 \\
1+\xi^2 \theta+\xi \theta^2 & 1+\xi^2 \theta & 1+\xi \theta^2 & 1-\xi \theta^2 & 1-\xi^2 \theta & 1-\xi^2 \theta-\xi \theta^2
\end{pmatrix}
$$

with $|A| \leq 3\sqrt{2}$ and $\|Ax\| \geq \sqrt{3}$ for every $x \in \mathbb{Z}_3^6$. 
Upper bound on $\|Ax\|$

While these results show the existence of matrices $A$ such that $\|Ax\|$ is bounded away from $\mathbf{0}$ on sparse vectors, it is also clear that for any $m \times d$ matrix $A$ there exist sparse vectors with $\|Ax\|$ not too large: for instance, if $x \in \mathbb{Z}^d$ is a standard basis vector, then

$$\|Ax\| \leq \sqrt{m} |A|. \quad (4)$$
Upper bound on $\|Ax\|$

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$$\|Ax\| \leq \sqrt{m} \ |A|.$$  \hspace{1cm} (4)

We prove a determinantal upper bound on $\|Ax\|$ in the spirit of Minkowski’s Geometry of Numbers, which is often better.

**Theorem 10 (F., Needell, Sudakov – 2017)**

Let $A$ be an $m \times d$ real matrix of rank $m \leq d$, and let $A'$ be the $d \times m$ real matrix so that $AA'$ is the $m \times m$ identity matrix. There exists a nonzero point $x \in \mathbb{Z}_m^d$ such that

$$\|Ax\| \leq \sqrt{m} \left| \det \left( (A')^\top A' \right) \right|^{-1/2m}.$$ \hspace{1cm} (5)
Upper bound examples

Example 2

Let \( d = 5 \), \( m = 3 \), and let

\[
A = \begin{pmatrix}
15 & 15 & 4 & 13 & 15 \\
2 & -1 & -15 & 2 & -13 \\
-13 & 2 & 1 & -15 & 4
\end{pmatrix},
\]

then

\[
A' = \begin{pmatrix}
3392/3905 & 23/355 & 3021/3905 \\
-1949/2130 & 3/710 & -1697/2130 \\
-6409/9372 & -19/284 & -5647/9372 \\
-6407/9372 & -17/284 & -6353/9372 \\
13869/15620 & 1/1420 & 12047/15620
\end{pmatrix},
\]

and so the bound of (5) is 8.375..., which is better than 25.980..., the bound given by (4).
**Upper bound examples**

**Example 3**

Let $d = 6$, $m = 3$, and let

$$A = \begin{pmatrix}
50000 & 20 & 40 & 3 & -50000 & 30 \\
-1 & -50000 & 20 & 40 & 4 & -50000 \\
-50000 & -1 & -50000 & -50000 & 20 & 40
\end{pmatrix},$$

then the bound of (5) is 7651.170... and (4) is 86602.540...
Upper bound examples

Example 3

Let $d = 6$, $m = 3$, and let

$$A = \begin{pmatrix} 50000 & 20 & 40 & 3 & -50000 & 30 \\ -1 & -50000 & 20 & 40 & 4 & -50000 \\ -50000 & -1 & -50000 & -50000 & 20 & 40 \end{pmatrix},$$

then the bound of (5) is 7651.170... and (4) is 86602.540...

Let $d = 8$, $m = 4$, and let

$$A = \begin{pmatrix} 6 & 13 & 13 & 11 & 6 & 12 & 11 & 10 \\ 7 & 12 & 6 & 13 & 7 & 11 & 11 & 9 \\ 8 & 11 & 12 & 9 & 12 & 12 & 11 \\ 13 & 10 & 7 & 8 & 13 & 13 & 13 & 13 \end{pmatrix},$$

then the bound of (5) is 2.412... and (4) is 26.
Sparse Geometry of Numbers

Let us now sketch the strategy of proof of Theorem 10. For this, we develop sparse analogues of some classical theorems in Minkowski’s Geometry of Numbers. We start with a result of J. D. Vaaler (1979).
Let us now sketch the strategy of proof of Theorem 10. For this, we develop sparse analogues of some classical theorems in Minkowski’s Geometry of Numbers. We start with a result of J. D. Vaaler (1979).

**Lemma 11 (Vaaler’s Cube Slicing Inequality)**

Let $C^d(1)$ be a cube of sidelength 1 centered at the origin in $\mathbb{R}^d$, i.e.

$$C^d(1) = \{ x \in \mathbb{R}^d : |x_i| \leq 1/2 \ \forall \ 1 \leq i \leq d \}.$$  

Let $V$ be an $m$-dimensional subspace of $\mathbb{R}^d$, $m \leq d$. Then the $m$-dimensional volume of the section $C_d(1) \cap V$ is

$$\text{Vol}_m(C_d(1) \cap V) \geq 1.$$
Sparse Geometry of Numbers

We can use Vaaler’s lemma to prove a sparse version of Minkowski’s Convex Body Theorem for parallelepipeds.

Proposition 12 (F., Needell, Sudakov – 2017)

Let $m \leq d$ be positive integers. Let $A \in \text{GL}_d(\mathbb{R})$, and let $P_A = AC^d(1)$. Assume that for some $I \subset [d]$ with $|I| = m$,

$$\sqrt{|\det(A_I^T A_I)|} \geq 2^m. \quad (6)$$

Then $P_A$ contains a nonzero point of $\mathbb{Z}_m^d$.

Here $[d] := \{1, \ldots, d\}$ and $A_I$ is the $d \times m$ submatrix of $A$ with columns indexed by elements of $I$. 
Sparse Geometry of Numbers

This implies a sparse version of Minkowski’s Linear Forms Theorem.

**Theorem 13 (F., Needell, Sudakov – 2017)**

Let $m \leq d$ be positive integers and $B \in \text{GL}_d(\mathbb{R})$. For each $1 \leq i \leq d$, let $L_i(X_1, \ldots, X_d) = \sum_{j=1}^{d} B_{ij} X_j$ be the linear form with entries of the $i$-th row of $B$ for its coefficients. Let $c_1, \ldots, c_d$ be positive real numbers such that for some $I = \{1 \leq j_1 < \cdots < j_m \leq d\} \subset [d]$,

$$c_{j_1} \cdots c_{j_m} \geq \left| \det \left( (B^{-1})_I^\top (B^{-1})_I \right) \right|^{-1/2}.$$  \hspace{1cm} (7)

Then there exists a nonzero point $x \in \mathbb{Z}_m^d$ such that

$$|L_{j_i}(x)| \leq c_{j_i} \quad \forall \ 1 \leq i \leq m.$$ \hspace{1cm} (8)
Sparse Geometry of Numbers

Taking \( I = \{1, \ldots, m\} \) and

\[
c_i = \left| \det \left( (B^{-1})_I^\top (B^{-1})_I \right) \right|^{-1/2m}
\]

for each \( 1 \leq i \leq m \) in Theorem 13 implies –

Corollary 14 (F., Needell, Sudakov – 2017)

Let \( A \) be an \( m \times d \) real matrix of rank \( m \leq d \). Let \( B \in \text{GL}_d(\mathbb{R}) \) be a matrix whose first \( m \) rows are the rows of \( A \). Let \( I = \{1, \ldots, m\} \). Then there exists a nonzero point \( x \in \mathbb{Z}_m^d \) such that

\[
\|Ax\| \leq \sqrt{m} \left| \det \left( (B^{-1})_I^\top (B^{-1})_I \right) \right|^{-1/2m}.
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Sparse Geometry of Numbers

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\( \) Let \( A \) be an \( m \times d \) real matrix of rank \( m \leq d \). Let \( B \in \text{GL}_d(\mathbb{R}) \) be a matrix whose first \( m \) rows are the rows of \( A \). Let \( I = \{1, \ldots, m\} \). Then there exists a nonzero point \( x \in \mathbb{Z}_m^d \) such that

\[
\|Ax\| \leq \sqrt{m} \left| \det \left( (B^{-1})_I^\top (B^{-1})_I \right) \right|^{-1/2m}.
\]

Theorem 10 follows immediately.
References


Preprints are available at:

http://math.cmc.edu/lenny/research.html
References


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Thank you!