# GRAPH PRODUCTS, FOURIER ANALYSIS AND SPECTRAL TECHNIQUES 

N. Alon, I. Dinur, E. Friedgut and B. Sudakov


#### Abstract

We consider powers of regular graphs defined by the weak graph product and give a characterization of maximum-size independent sets for a wide family of base graphs which includes, among others, complete graphs, line graphs of regular graphs which contain a perfect matching and Kneser graphs. In many cases this also characterizes the optimal colorings of these products.

We show that the independent sets induced by the base graph are the only maximum-size independent sets. Furthermore we give a qualitative stability statement: any independent set of size close to the maximum is close to some independent set of maximum size.

Our approach is based on Fourier analysis on Abelian groups and on Spectral Techniques. To this end we develop some basic lemmas regarding the Fourier transform of functions on $\{0, \ldots, r-1\}^{n}$, generalizing some useful results from the $\{0,1\}^{n}$ case.


## 1 Introduction

Consider the following combinatorial problem:
Assume that at a given road junction there are $n$ three-position switches that control the red-amber-green position of the traffic light. You are told that whenever you change the position of all the switches then the color of the light changes. Prove that in fact the light is controlled by only one of the switches.

The above problem is a special case of the problem we wish to tackle in this paper, characterizing the optimal colorings and maximal independent

[^0]sets of products of regular graphs. The configuration space of the switches described above can be modeled by the $n$-fold product of $K_{3}$. Let us begin by defining the weak graph product of two graphs.

The weak product of $G$ and $H$, denoted by $G \times H$ is defined as follows: the vertex set of $G \times H$ is the Cartesian product of the vertex sets of $G$ and $H$. Two vertices $\left(g_{1}, h_{1}\right)$ and $\left(g_{2}, h_{2}\right)$ are adjacent in $G \times H$ if $g_{1} g_{2}$ is an edge of $G$ and $h_{1} h_{2}$ is an edge of $H$. The "times" symbol, $\times$, is supposed to be reminiscent of the weak product of two edges: $\mid \times-=\times$. In this paper "graph product" will always mean the weak product.

In the first part of the paper we consider the interesting special case of the product of complete graphs on $r>2$ vertices,

$$
G=K_{r}^{n}=\times_{j=1}^{n} K_{r} .
$$

We then discuss a more general setting, considering other $r$-regular graphs as well.

When $G=K_{r}^{n}$, we identify the vertices of $G$ in the obvious way with the elements of $\mathbb{Z}_{r}^{n}$. Recalling the definition of the product, two vertices are adjacent in $G$ iff the corresponding vectors differ in every coordinate. Clearly one can color $G$ with $r$ colors by choosing a coordinate $i$ and coloring every vertex according to its $i$ th coordinate. The following theorem asserts that if $r>2$ then these are the only $r$-colorings. Here, and in what follows, we denote by $|H|$ the number of vertices of a graph $H$.
Theorem 1.1. Let $G=K_{r}^{n}$, and assume $r \geq 3$. Let $I$ be an independent set with $|I|=|G| / r$. Then there exists a coordinate $i \in\{1, \ldots, n\}$ and $k \in\{0, \ldots, r-1\}$ such that

$$
I=\left\{v: v_{i}=k\right\} .
$$

Consequently, the only colorings of $G$ by $r$ colors are those induced by colorings of one of the factors $K_{r}$.

Greenwell and Lovász [GL] proved the above theorem (and actually, a somewhat stronger statement) more than a quarter of a century ago. See also [Mü] for a similar result. The novelty in this paper is the proof we supply that uses Fourier analysis on the group $\mathbb{Z}_{r}^{n}$. Our approach also allows us to deduce a stability version of the above theorem.
Theorem 1.2. For every $r \geq 3$ there exists a constant $M=M(r)$ such that for any $\epsilon>0$ the following is true. Let $G=K_{r}^{n}$. Let $J$ be an independent set such that $\frac{|J|}{|G|}=\frac{1}{r}-\epsilon$. Then there exists an independent set $I$ with $\frac{|I|}{|G|}=\frac{1}{r}$ such that $\frac{|J \triangle I|}{|G|}<M \epsilon$.

Here " $\triangle$ " denotes the symmetric difference. What the above theorem tells us is (in conjunction with Theorem 1.1) that any independent set that is close to being of maximum-size is close to being determined by one coordinate. We do not know of any purely combinatorial proof of this result.

The results in both theorems above can be extended to other base graphs. Let $\alpha(G)$ denote the maximum possible size of an independent set in a graph $G$. The following observation determines $\alpha\left(H^{n}\right)$ for any vertex transitive base graph $H$, in terms of $\alpha(H)$ and $|H|$.
Proposition 1.3. For any vertex transitive graph $H$ and for any integer $n \geq 1$, if $G=H^{n}$ then

$$
\frac{\alpha(G)}{|G|}=\frac{\alpha(H)}{|H|} .
$$

After the simple proof of this proposition (some special cases of which are proved in $[\mathrm{BrNR}]$ ), we will provide some examples showing that the above equality does not necessarily hold without the transitivity assumption.

The relevance of graph eigenvalues to independent sets in graphs is well known and can be traced back to the old result that the independence number of any regular graph $H$ on $r$ vertices in which the eigenvalues of the adjacency matrix are $\mu_{1} \geq \mu_{2} \geq \cdots \geq \mu_{r}$, is at most $-r \mu_{r} /\left(\mu_{1}-\mu_{r}\right)$. A proof of this fact, as well as of the related results on the connection between the Shannon capacity of a graph and its eigenvalues, can be found in [Lo]. This bound is tight for many graphs $H$ including, for example, complete graphs and the Petersen graph. It turns out that the results in Theorem 1.1 and in Theorem 1.2 can be extended to any connected non-bipartite regular base graph $H$ for which the above bound is tight.
Theorem 1.4. Let $H$ be a connected $d$-regular graph on $r$ vertices and let $d=\mu_{1} \geq \mu_{2} \geq \cdots \geq \mu_{r}$ be its eigenvalues. If

$$
\begin{equation*}
\frac{\alpha(H)}{r}=\frac{-\mu_{r}}{d-\mu_{r}} \tag{1}
\end{equation*}
$$

then for every integer $n \geq 1$,

$$
\frac{\alpha\left(H^{n}\right)}{r^{n}}=\frac{-\mu_{r}}{d-\mu_{r}} .
$$

Moreover, if $H$ is also non-bipartite, and if $I$ is an independent set of size $\frac{-\mu_{r}}{d-\mu_{r}} r^{n}$ in $G=H^{n}$, then there exists a coordinate $i \in\{1,2, \ldots, n\}$ and a maximum independent set $J$ in $H$, such that

$$
I=\left\{v \in V(H)^{n}: v_{i} \in J\right\} .
$$

Remark. Note that for any $H$ and $n, \chi\left(H^{n}\right)=\chi(H)$. If $H$ satisfies the conditions of the last theorem and if, in addition, $\chi(H)=r / \alpha(H)$ then every optimal coloring of $H^{n}$ is induced by a coloring of one of the multiplicands, since it is a partition of $H^{n}$ into maximum-size independent sets. Such a partition can only be consistent if each color class is induced by the same coordinate. The assumption $\chi(H)=r / \alpha(H)$ holds for many of the interesting classes of graphs to which Theorem 1.4 applies, see subsection 5.1.
Theorem 1.5. Let $H$ be a d-regular, connected, non-bipartite graph on $r$ vertices, let $d=\mu_{1} \geq \mu_{2} \geq \cdots \geq \mu_{r}$ be its eigenvalues and suppose its independence number satisfies (1). Then, there exists a constant $M=$ $M(H)$ such that for any $\epsilon>0$ the following holds. Let $G=H^{n}$ and let $I$ be an independent set such that $\frac{|I|}{|G|}=\frac{\alpha(H)}{|H|}-\epsilon$. Then there exists an independent set $I^{\prime}$ with $\frac{\left|I^{\prime}\right|}{|G|}=\frac{\alpha(H)}{|H|}$ such that $\frac{\left|I^{\prime} \Delta I\right|}{|G|}<M \epsilon$.

The rest of this paper is organized as follows. Section 2 is devoted to Fourier analysis of $\mathbb{Z}_{r}^{n}$, where we set up some of the necessary tools. We continue in section 3 with the proofs of several lemmas regarding Boolean functions on $\mathbb{Z}_{r}^{n}$. In section 4 we present the proofs of Theorem 1.1 and Theorem 1.2, and suggest a conjecture that strengthens Theorem 1.2 and offers a characterization of all independent sets that occupy a constant proportion of $G=K_{r}^{n}$. We also briefly consider the case of infinite dimensional products. Section 5 contains the proof of Proposition 1.3 and the proofs of Theorems 1.4 and 1.5, together with some related examples. These proofs are based on spectral techniques. Although the proofs of Theorems 1.1 and 1.2 can be obtained as special cases, as the eigenvectors of Cayley graphs of Abelian groups (and in particular, of complete graphs) are the characters of the corresponding groups, we believe that presenting them separately is instructive. It is interesting to note that the approach in section 5 suggests that many tools of Harmonic Analysis hold for general tensor products of some fixed orthogonal set of vectors, and not only for characters. Although this fact is not difficult and its proof essentially amounts to a change of basis, it seems very powerful and may well lead to additional interesting consequences. The final section 6 contains some concluding remarks.

## 2 Fourier Analysis on $\mathbb{Z}_{r}^{n}$

In this section we describe all the necessary background concerning discrete Fourier analysis. In addition, we present two extensions of known
results for functions $f:\{0,1\}^{n} \rightarrow\{0,1\}$ to the equivalent results for functions $f: \mathbb{Z}_{r}^{n} \rightarrow\{0,1\}$. The first, Lemma 2.2, is a generalization of the Bonami-Beckner hyper-contractive estimate, and the second, Lemma 2.4, is a generalization of a result proved in [FKN] regarding functions whose Fourier representation is concentrated on subsets of size $\leq 1$. We believe that these generalizations may find further uses. The proofs of these lemmas are presented in section 3.
2.1 Background and lemmas. Let $r>2$ and $G=\{0, \ldots, r-1\}^{n}$ $=\mathbb{Z}_{r}^{n}$. Here is some notation we shall use: Let $\overline{0}=(0,0, \ldots, 0), \overline{1}=$ $(1,1, \ldots, 1)$, and $e_{i}=(0, \ldots, 1, \ldots, 0)$, the unit vector with 1 at the $i$ th coordinate. For $S \in G$ let

$$
|S|=\left|\left\{i: S_{i} \neq 0\right\}\right| .
$$

We will occasionally refer to the "levels" of $\{0, \ldots, r-1\}^{n}$ : The 0-level contains only $\overline{0}$, the $k$ th level all $S$ such that $|S|=k$.

We think of $G$ both as an Abelian group and as a probability space endowed with the uniform (product) measure $\mu$. For any function $f: G \rightarrow \mathbb{C}$, let

$$
\int_{G} f(S) d S \stackrel{\text { def }}{=} \frac{1}{|G|} \sum_{S \in G} f(S)
$$

For any such function $f$ we define the $p$-norm of $f$ :

$$
\|f\|_{p}=\left(\int|f(S)|^{p}\right)^{1 / p}
$$

The inner product between two functions $f$ and $g$ is

$$
\langle f, g\rangle=\int_{G} f(S) \overline{g(S)} d S=\frac{1}{|G|} \sum_{S \in G} f(S) \overline{g(S)} .
$$

We now wish to consider the dual group to $G$, the group of characters, whose elements form an orthonormal basis for the space of complex-valued functions on $G$. For any $S \in G$ let $u_{S}: G \rightarrow \mathbb{C}$ be defined by

$$
u_{S}(T)=e^{2 \pi i\langle S, T\rangle / r}
$$

where $\langle S, T\rangle=\sum_{i=1}^{n} S_{i} T_{i} \bmod r$. It is easy to verify the following properties of the characters:

- They indeed form a group, and the mapping $S \mapsto u_{S}$ is a homomorphism from $G$ to the multiplicative group of roots of unity:

$$
u_{S} u_{T}=u_{S+T}
$$

and

$$
u_{-T}=u_{T}^{-1}=\overline{u_{T}} .
$$

- Since the sum of the roots of unity is 0 we have that for $S \neq \overline{0}$, $\sum_{T \in G} u_{S}(T)=0$. Of course $u_{\overline{0}} \equiv 1$, hence $\frac{1}{|G|} \sum_{T \in G} u_{\overline{0}}(T)=1$.
- It follows that the characters form an orthonormal basis:

$$
\left\langle u_{S}, u_{T}\right\rangle=\delta_{S, T} .
$$

- Every function $f: G \rightarrow \mathbb{C}$ has a unique expansion of the form $f=$ $\sum \widehat{f}(S) u_{S}$, where $\widehat{f}(S)=\left\langle f, u_{S}\right\rangle=\frac{1}{\mid G} \sum_{T \in G} f(T) \cdot \overline{u_{S}(T)}$.
- From orthonormality we get that, for any two functions $f, g$,

$$
\langle f, g\rangle=\sum \widehat{f}(S) \overline{\hat{g}(S)},
$$

and in particular we get Parseval's identity,

$$
\|f\|_{2}^{2}=\langle f, f\rangle=\sum|\widehat{f}(S)|^{2} .
$$

The usage of Fourier analysis in the study of Boolean functions and various questions in combinatorics and computer science has seen much success in the last fifteen years, see for example [AKRS], [BenKS], [Bou2], [F],[FKN], [KKL], [KiS], [LiMN], [T]. All the above examples deal with problems set in $\{0,1\}^{n}$. In contrast there are relatively few combinatorial papers dealing with $\mathbb{Z}_{r}^{n}$ that use Fourier analysis. Some examples of the latter can be found in $[\mathrm{BS}],[\mathrm{ALM}],[\mathrm{M}]$.

A central tool that appears in almost every one of the aforementioned works that deal with $\{0,1\}^{n}$ is a hyper-contractive estimate due, independently, to Beckner and Bonami, $[\mathrm{Be}]$, $[\mathrm{Bo}]$. One version of the BonamiBeckner inequality that is often used is the following (see also [Bou1]):
Lemma 2.1 (Beckner (1975), Bonami (1970)). Let $f:\{0,1\}^{n} \rightarrow \mathbb{R}$ be a function that is a linear combination of $\left\{u_{T}:|T| \leq k\right\}$. Let $p>2$. Then

$$
\|f\|_{p} \leq(\sqrt{p-1})^{k}\|f\|_{2}
$$

Bonami and Beckner actually proved the dual version of the above lemma, which is equivalent to it. Here we will need a more general version dealing with functions on $\mathbb{Z}_{r}^{n}$. On the other hand, as we do not need optimal constants or a statement regarding all $p$-norms, we will make do with the following statement. Recall that $|T|=\left|\left\{i: T_{i} \neq 0\right\}\right|$.
Lemma 2.2. For every $r \geq 2$ there exists $C>0$ such that the following holds. Let $G=\mathbb{Z}_{r}^{n}$ and $f: G \rightarrow \mathbb{C}$ be a function whose Fourier transform is concentrated on the first $k+1$ levels, that is, $f$ is a linear combination of $\left\{u_{T}:|T| \leq k\right\}$. Then

$$
\|f\|_{4} \leq C^{k}\|f\|_{2}
$$

We will use this inequality in the proof of Lemma 2.4 below.

The following two lemmas give conditions on the Fourier transform of Boolean functions $f$ under which $f$ may be determined or approximated by a function depending on only one coordinate, a "dictator".
Lemma 2.3. Let $f: \mathbb{Z}_{r}^{n} \rightarrow\{0,1\}$ be such that all its Fourier transform is concentrated on the first two levels,

$$
|S|>1 \Rightarrow \widehat{f}(S)=0 .
$$

Then $f$ is either constant or depends on precisely one coordinate.
The following lemma is a stability version of the above which ultimately enables us to deduce Theorem 1.2 and Theorem 1.5. The special case of this lemma with $G=\{0,1\}^{n}$ is stated and proven in [FKN].
Lemma 2.4. Let $r \geq 2$ and let $C=C(r)$ be the constant from Lemma 2.2. Let $K=2+32 C^{8}$. Then for any $\varepsilon>0$ the following holds: let $f: \mathbb{Z}_{r}^{n} \rightarrow\{0,1\}$ be a function such that $\operatorname{Pr}[f=1]=\alpha$ and furthermore assume that

$$
\sum_{|S|>1}|\widehat{f}(S)|^{2}=\varepsilon .
$$

Then there exists a function $g: \mathbb{Z}_{r}^{n} \rightarrow\{0,1\}$ which depends on at most one coordinate such that

$$
\|f-g\|_{2}^{2}<\frac{K}{\alpha-\alpha^{2}-\varepsilon} \varepsilon .
$$

Remark. Note that under the above assumptions $\varepsilon \leq \alpha-\alpha^{2}$, because $\alpha-\alpha^{2}-\varepsilon=\sum_{|S|=1}|\widehat{f}(S)|^{2}$.

## 3 Generalizing the Boolean Inequalities

In this section we prove several lemmas regarding Boolean functions on $\mathbb{Z}_{r}^{n}$. The first, Lemma 2.2, generalizes a hyper-contractive estimate due to Beckner and Bonami [Be], [Bo]. Lemma 2.4 then generalizes a result of Friedgut, Kalai and Naor [FKN] regarding functions whose Fourier representation is concentrated on the levels 0 and 1 . We begin by proving Lemma 2.2 .

Proof of Lemma 2.2. Following the original proofs of Beckner and Bonami we will prove the dual statement formalized in Claim 3.1 below, using induction on the dimension. For $0<\varepsilon<1$ and

$$
f(s)=\sum_{j=0}^{r-1} a_{j} e^{\frac{2 \pi i j s}{r}}:\{0, \ldots, r-1\} \rightarrow \mathbb{C}
$$

define

$$
T_{\delta} f(s)=a_{0}+\delta\left(\sum_{j=1}^{r-1} a_{j} e^{\frac{2 \pi i j s}{r}}\right)
$$

By abuse of notation we will also denote $T_{\delta}^{\otimes n}$ by $T_{\delta}$ : for any function $f:\{0, \ldots, r-1\}^{n} \rightarrow \mathbb{C}$ let $T_{\delta} f=\sum \widehat{f}(S) \delta^{|S|} u_{S}$.
Claim 3.1. For every $r \geq 2$ there exists $0<\delta<1$ such that, for every $f:\{0, \ldots, r-1\}^{n} \rightarrow \mathbb{C}$,

$$
\left\|T_{\delta} f\right\|_{4} \leq\|f\|_{2} .
$$

Remarks. This claim immediately implies Lemma 2.2 with $C=1 / \delta$. Indeed,

$$
\begin{aligned}
\|f\|_{4}=\left\|T_{\delta}\left(T_{\delta^{-1}} f\right)\right\|_{4} & \leq\left\|T_{\delta^{-1}} f\right\|_{2}=\left(\sum_{|S| \leq k} \delta^{-2|S|}|\widehat{f}(S)|^{2}\right)^{1 / 2} \\
& \leq(1 / \delta)^{k}\left(\sum_{|S| \leq k}|\widehat{f}(S)|^{2}\right)^{1 / 2}=(1 / \delta)^{k}\|f\|_{2} .
\end{aligned}
$$

Furthermore, we interpret Claim 3.1 as saying that $\left\|T_{\delta}\right\|=1$ as an operator from $L_{4}$ to $L_{2}$. As observed by Bonami and Beckner, for the class of operators defined by integration against a kernel (such as $T_{\delta}$ ) one has multiplicativity of operator norms. (This can be proven rather easily using Minkowski's inequality.) Therefore one can now use induction on the dimension, and it suffices to prove the one-dimensional statement.
Proof of the claim. Fix $r \geq 2$. Let $\omega=e^{2 \pi i / r}$ and let $f(s)=\sum_{j=0}^{r-1} a_{j} \omega^{s j}$. Clearly we may assume that $a_{0}=1$ and that $\left|a_{1}\right| \geq\left|a_{j}\right|$ for $j \geq 1$. Furthermore, since replacing $a_{i}$ by $\left|a_{i}\right|$ does not change the norms of neither $f$ nor $T_{\delta} f$, we may assume that the $a_{i}$ 's are non-negative reals.

Let $1>\delta>0$ and $T_{\delta} f(s)=\sum_{j=0}^{r-1} b_{j} \omega^{s j}$. Then $b_{0}=a_{0}=1$ and $b_{j}=\delta a_{j}$ for $j \geq 1$. A simple calculation gives

$$
\left[T_{\delta} f(s)\right]^{2}=\sum_{j}\left[\sum_{k} b_{k} b_{j-k}\right] \omega^{s j}
$$

We then have the following:

$$
\left\|T_{\delta} f\right\|_{4}^{4}=\left\|\left[T_{\delta} f\right]^{2}\right\|_{2}^{2}=\sum_{j}\left[\sum_{k} b_{k} b_{j-k}\right]^{2}
$$

whereas

$$
\|f\|_{2}^{4}=\left(\sum a_{j}^{2}\right)^{2}
$$

Denote $a_{1}=a$ and $b_{1}=b=\delta a$. Then

$$
1+a^{2}+a^{4} \leq\|f\|_{2}^{4},
$$

and, recalling that $b_{1}=b \geq b_{j}$ for $j \geq 1$,

$$
\left\|T_{\delta} f\right\|_{4}^{4} \leq 1+K\left[b^{2}+b^{3}+b^{4}\right] \leq 1+K \delta\left[a^{2}+a^{3}+a^{4}\right]
$$

where $K$ is some constant depending on $r$ (but not on $f$.) Hence what we need to show is that for sufficiently small $\delta$

$$
1+K \delta\left[a^{2}+a^{3}+a^{4}\right] \leq 1+a^{2}+a^{4}
$$

holds for all values of $a$, or equivalently, that

$$
(1-K \delta) a^{2}-K \delta a+(1-K \delta)>0
$$

The discriminant of this quadratic equation is

$$
4 K^{2} \delta^{2}-(1-K \delta)^{2}
$$

which is negative for sufficiently small values of $\delta$, which proves the assertion of the lemma.

We now turn to the proof of Lemma 2.3, which is a softened version of Lemma 2.4.
Proof. The function $f$ is of the form

$$
f(S)=a_{0}+\sum_{j=1}^{n} \sum_{k=1}^{r-1} a_{j, k} e^{\frac{2 \pi i S_{j} k}{r}}
$$

Since $f$ is Boolean we have $f=f^{2}$. From the uniqueness of the Fourier expansion we may compare the coefficients of $e^{2 \pi i S_{j_{1}} k / r} e^{2 \pi i S_{j_{2}} l / r}$ in $f$ and $f^{2}$ and deduce that if $j_{1} \neq j_{2}$ then

$$
a_{j_{1}, k} a_{j_{2}, l}=0
$$

Therefore there must exist one coordinate, say $j_{l}$, such that if $j \neq j_{l}$ then $a_{j, k}=0$ for all $k$. Hence $f$ is of the form

$$
f(S)=a_{0}+\sum_{k=1}^{r-1} a_{k} e^{\frac{2 \pi i S_{l} k}{r}}
$$

Finally we turn to the proof of Lemma 2.4. Our proof follows the lines of one of the proofs presented in [FKN].
Proof of Lemma 2.4. To begin we may assume that $\varepsilon \leq 1 / 4 C^{8}$ where $C$ is the constant from Lemma 2.2. If this assumption fails then letting $K=2+32 C^{8}$, makes finding a $K \varepsilon$-approximation of $f$ trivial. Let

$$
f_{S}=\sum_{|T| \leq 1} \widehat{f}(T) u_{T} \quad \text { and } \quad f_{L}=\sum_{|T|>1} \widehat{f}(T) u_{T}
$$

( $S$ and $L$ stand for small and large.)

If $f_{S}$ were Boolean, then the lemma would follow easily. We will show that the same is true for $f_{S}$ that is close to being Boolean. Let us consider the following function, $h$, that we use to measure, in a sense, how far $f_{S}$ is from being Boolean:

$$
h=f_{S}^{2}-f_{S}
$$

Note that if $f_{S}$ were Boolean then $h$ would be identically 0 .
We start by computing $\widehat{h}(S)$ on the second level. Let $i \neq j$. Then

$$
\begin{equation*}
\widehat{h}\left(k_{1} e_{i}+k_{2} e_{j}\right)=\widehat{f_{S}^{2}}\left(k_{1} e_{i}+k_{2} e_{j}\right)-\widehat{f_{S}}\left(k_{1} e_{i}+k_{2} e_{j}\right)=2 \widehat{f}\left(k_{1} e_{i}\right) \widehat{f}\left(k_{2} e_{j}\right) \tag{2}
\end{equation*}
$$

where we have used the fact that $\widehat{f_{S}}$ is 0 on the second level. The following lemma shows that the $L_{2}$ norm of $h$ is small.
Lemma 3.2. Let $C$ be the constant from Lemma 2.2 and let $\lambda=32 C^{8}$. Then

$$
E\left(|h|^{2}\right) \leq \lambda \varepsilon
$$

Assuming Lemma 3.2 it is now easy to prove that the dictatorapproximation of $f$ guaranteed by Lemma 2.4 can be obtained from $f$ by "truncating and rounding":
Corollary 3.3. There exists $1 \leq j \leq n$ such that the for the function $g^{\prime}$, defined as follows,

$$
\begin{equation*}
g^{\prime}(x)=\widehat{f}(\overline{0})+\sum_{k=1}^{r-1} \widehat{f}\left(k e_{j}\right) e^{\frac{2 \pi i k x_{j}}{r}} \tag{3}
\end{equation*}
$$

it is true that $\left\|g^{\prime}-f\right\|_{2}^{2} \leq \varepsilon\left(1+\frac{\lambda}{2\left(\alpha-\alpha^{2}-\varepsilon\right)}\right)$.
This corollary now implies Lemma 2.4. It is not hard to see that if $g(y)$ is the function obtained by rounding $g^{\prime}$ to 0 or 1 according to the majority value of $f(x)$ over all $\left\{x \mid x_{j}=y_{j}\right\}$ then $\|g-f\|_{2}^{2} \leq 2\left\|g^{\prime}-f\right\|_{2}^{2} \leq \frac{K}{\alpha-\alpha^{2}-\varepsilon} \varepsilon$.
Proof of the corollary. Denote $a_{i}=\sum_{k=1}^{r-1}\left|\widehat{f}\left(k e_{i}\right)\right|^{2}$. Since $E\left(|h|^{2}\right)=$ $\sum|\widehat{h}(T)|^{2}$ we have, by summing only over $T$ of the form $T=k_{1} e_{i}+k_{2} e_{j}$ and using (2)

$$
\begin{equation*}
\sum_{i<j} a_{i} a_{j}=\sum_{i<j}\left[\sum_{k=1}^{r-1}\left|\widehat{f}\left(k e_{i}\right)\right|^{2}\right]\left[\sum_{k=1}^{r-1}\left|\widehat{f}\left(k e_{j}\right)\right|^{2}\right] \leq \frac{E\left(|h|^{2}\right)}{4} \leq \frac{\lambda}{4} \varepsilon \tag{4}
\end{equation*}
$$

We also know that $\widehat{f}(\overline{0})=\alpha$ and that $\sum|\widehat{f}(T)|^{2}=\alpha$. This means that

$$
\sum_{i} a_{i}=\alpha-\alpha^{2}-\epsilon
$$

Let $j$ be the index for which $a_{j}$ is maximal.

$$
\left(\sum_{i} a_{i}\right)^{2} \leq \sum_{i} a_{i}^{2}+\frac{\lambda}{2} \varepsilon \leq a_{j} \cdot \sum_{i} a_{i}+\frac{\lambda}{2} \varepsilon
$$

which means, dividing by $\sum_{i} a_{i}=\alpha-\alpha^{2}-\varepsilon$, that $a_{j} \geq \alpha-\alpha^{2}-\varepsilon\left(1+\frac{\lambda}{2\left(\alpha-\alpha^{2}-\varepsilon\right)}\right)$. This means that $\left\|g^{\prime}\right\|_{2}^{2}=a_{j}+\alpha^{2} \geq \alpha-\varepsilon\left(1+\frac{\lambda}{2\left(\alpha-\alpha^{2}-\varepsilon\right)}\right)$ hence most of $f$ 's $L_{2}$ norm comes from the Fourier transform of $g^{\prime}$, implying the corollary.
Proof of Lemma 3.2. We first show that $h$ is 'mostly' smaller than $O(\sqrt{\varepsilon})$, and then use a concentration of measure argument to show that it cannot be much larger on the rest of the domain.

Recall that $f_{S}=f-f_{L}$ and that $f^{2}-f=0$, so

$$
h=\left(f_{S}\right)^{2}-f_{S}=\left(f-f_{L}\right)^{2}-\left(f-f_{L}\right)=f_{L}^{2}+f_{L}(1-2 f)
$$

Let $k=2 C^{4}$, and set $Z=\left\{x \in \mathbb{Z}_{r}^{n}| | f_{L}(x) \mid \leq k \sqrt{\varepsilon}\right\}$. Since $\left\|f_{L}\right\|_{2}^{2} \leq \varepsilon$, a Markov argument implies $\operatorname{Pr}_{x}[Z] \geq 1-1 / k^{2}$. Additionally, for every $x \in Z,|h(x)| \leq 2\left|f_{L}(x)\right| \leq 2 k \sqrt{\varepsilon}$. (Here we have used our assumption that $\varepsilon<1 / k^{2}=C^{-8} / 4$, see the first sentence of the proof.)

How can $E\left(|h|^{2}\right)$ be much larger than $\varepsilon$ ? Only if most of the contribution to this expectation comes from values of $x$ outside $Z$, where $h$ obtains values much larger than $\sqrt{\varepsilon}$. We next use the hyper-contractive estimate (Lemma 2.2) to show that this cannot happen; intuitively because having only low Fourier frequencies means $|h|$ is fairly 'concentrated' around its expectation.

For convenience of notation, let $X=E\left(|h|^{2}\right)$ and $Y=E\left(|h|^{4}\right)$, and denote $p=\operatorname{Pr}_{x}[x \notin Z] \leq 1 / k^{2}$. For every function $h$, it is true that $X \leq \sqrt{Y}$. Since the Fourier support of $h$ is only on the first two levels, Lemma 2.2 implies that $\sqrt{Y} \leq C^{4} X$. Thus,

$$
\begin{aligned}
X=E\left(|h|^{2}\right) & =(1-p) \cdot E\left(|h(x)|^{2} \mid x \in Z\right)+p \cdot E\left(|h(x)|^{2} \mid x \notin Z\right) \\
& \leq(1-p) 4 k^{2} \varepsilon+p \sqrt{E\left(|h(x)|^{4} \mid x \notin Z\right)} \\
& \leq 4 k^{2} \varepsilon+p \sqrt{\frac{Y}{p}} \\
& \leq 4 k^{2} \varepsilon+\sqrt{p} C^{4} X \leq 4 k^{2} \varepsilon+\frac{1}{2} X
\end{aligned}
$$

Thus, $X \leq 8 k^{2} \varepsilon=32 C^{8} \varepsilon$.

## 4 Independent Sets in $K_{r}^{\boldsymbol{n}}$

In this section we prove Theorems 1.1 and 1.2. Let $G=K_{r}^{n}$. Identify the vertices of $G$ with the elements of $\mathbb{Z}_{r}^{n}$. Let $D=\left(\mathbb{Z}_{r} \backslash 0\right)^{n}$. Recall that two
vertices $u, v \in \mathbb{Z}_{r}^{n}$ are adjacent iff $u-v \in D . G$ is thus the Cayley graph of $\mathbb{Z}_{r}^{n}$ with the elements of $D$ as generators, and letting $d=|D|=(r-1)^{n}$ we have that the graph $G$ is $d$ regular. Aside from its group structure we will also think of $\mathbb{Z}_{r}^{n}$ as a measure space with the uniform (product) measure, which we shall denote by $\mu$.

For any $i$ between 1 and $n$ and any $0 \leq k \leq r-1$, the set $V_{i}^{k}=\left\{v \mid v_{i}=k\right\}$ is an example of a large independent set, $\mu\left(V_{i}^{k}\right)=1 / r$. (Although the measure notation is unusual in the combinatorial setting, it will blend in well with our analytical methods in the sequel.) We first observe that this is the maximum size that an independent set can attain.
Claim 4.1. If $I$ is an independent set in $G$ then $\mu(I) \leq 1 / r$.
We now prove the claim by a simple combinatorial argument. However, it is interesting to note that the claim will also follow easily from the analytical proof of Theorem 1.1 that we present below.
Proof of the claim. We will partition the graph into disjoint cliques of size $r$. The claim will then follow from the fact that $I$ can intersect every clique in at most one vertex.

Let $\bar{i}=(i, i, \ldots, i) \in \mathbb{Z}_{r}^{n}$, and let $K=\{\overline{0}, \overline{1}, \overline{2}, \ldots, \overline{(r-1)}\}$. The claim is proven by observing that $v+K$ is a clique for any $v \in \mathbb{Z}_{r}^{n}$ and that

$$
\mathbb{Z}_{r}^{n}=\bigcup_{\left(v_{2}, \ldots, v_{n}\right) \in \mathbb{Z}_{r}^{n-1}}\left(\left(0, v_{2}, \ldots, v_{n}\right)+K\right)
$$

(Note that this is the partition of $\mathbb{Z}_{r}^{n}$ into cosets of the subgroup K.) $\quad \square$
Proposition 1.3, proved in section 5, generalizes the above claim for arbitrary regular vertex transitive graphs.

Theorems 1.1 and 1.2 will follow from the next lemma, that is inspired by observations of Håstad [Ha], and is also a special case of the basic spectral approach described in section 5.
Lemma 4.2. Let $I \subset \mathbb{Z}_{r}^{n}$ be an independent set in $G$, and let $f: \mathbb{Z}_{r}^{n} \rightarrow\{0,1\}$ be its indicator function, i.e. $f(x)=1$ iff $x \in I$. Then,

$$
\sum_{S}|\widehat{f}(S)|^{2}\left(\frac{-1}{r-1}\right)^{|S|}=0
$$

Proof. For any $\tau \in D$ define $f_{\tau}(x)=f(x+\tau)$. Since $I$ is an independent set, always $f_{\tau}(x) \cdot f(x)=0$, so $\left\langle f, f_{\tau}\right\rangle=0$. Define

$$
A(f)=\frac{1}{d} \sum_{\tau \in D} f_{\tau} .
$$

$A$ is the averaging operator, replacing $f(x)$ by the average of $f$ on the neighbors of $x$. The above implies that (and indeed is equivalent to)

$$
\begin{equation*}
\langle f, A(f)\rangle=0 . \tag{5}
\end{equation*}
$$

Let us now compute the Fourier transform of $A(f)$ in terms of the Fourier transform of $f$. Recall that for $S \in \mathbb{Z}_{r}^{n}$ we denote $|S|=\left|\left\{i \mid S_{i} \neq 0\right\}\right|$, Claim 4.3.

$$
\widehat{A(f)}(S)=\widehat{f}(S)\left(\frac{-1}{r-1}\right)^{|S|}
$$

Proof. We first claim that $\widehat{f}_{\tau}(S)=\widehat{f}(S) u_{\tau}(S)$. Indeed,

$$
\begin{aligned}
\widehat{f_{\tau}}(S)=\int f(x+\tau) \overline{u_{S}(x)} d x & =\int f(x) \overline{u_{S}(x-\tau)} d x \\
& =\int f(x) \overline{u_{S}(x)} u_{S}(\tau) d x=\widehat{f}(S) u_{\tau}(S)
\end{aligned}
$$

Recall that $D=\{1, \ldots, r-1\}^{n}, d=|D|$, and denote $\omega=e^{2 \pi i / r}$. We have

$$
\begin{aligned}
\widehat{A(f)}(S) & =\frac{1}{d} \sum_{\tau \in D} \widehat{f_{\tau}}(S)=\frac{1}{d} \widehat{f}(S) \sum_{\tau \in D} u_{\tau}(S) \\
& =\frac{1}{d} \widehat{f}(S) \prod_{j=1}^{n} \sum_{k=1}^{r-1} \omega^{k S_{j}} \\
& =\frac{1}{d} \widehat{f}(S) \prod_{j: S_{j}=0}(r-1) \prod_{j: S_{j} \neq 0}(-1)=\widehat{f}(S)\left(\frac{-1}{r-1}\right)^{|S|}
\end{aligned}
$$

where in the last equality we have used $d=(r-1)^{n}$. The proof of the lemma is now completed by orthogonality, writing

$$
0=\langle f, A(f)\rangle=\sum_{S} \widehat{f}(S) \overline{\widehat{A(f)}(S)}=\sum_{S}|\widehat{f}(S)|^{2}\left(\frac{-1}{r-1}\right)^{|S|}
$$

Proof of Theorems 1.1 and 1.2. Denote $\alpha=\mu(I)$. Now, from Parseval we have

$$
\sum|\widehat{f}(S)|^{2}=\alpha
$$

Also,

$$
\widehat{f}(\overline{0})=\int f(x) d x=\alpha .
$$

Hence

$$
\begin{equation*}
\sum_{S \neq \overline{0}}|\widehat{f}(S)|^{2}=\alpha-\alpha^{2} \tag{6}
\end{equation*}
$$

and subtracting $|\widehat{f}(0)|^{2}=\alpha^{2}$ from the identity of Lemma 4.2,

$$
\begin{equation*}
\sum_{S \neq \overline{0}}|\widehat{f}(S)|^{2}\left(\frac{-1}{r-1}\right)^{|S|}=-\alpha^{2} \tag{7}
\end{equation*}
$$

Comparing (6) and (7) now yields valuable information on the Fourier transform of $f$. It is convenient to interpret this information via the following probability distribution: Let $T$ be a random variable taking values in $G \backslash \overline{0}$ with

$$
\operatorname{Pr}[T=S]=\frac{|\widehat{f}(S)|^{2}}{\alpha-\alpha^{2}} .
$$

Let

$$
X=X(T)=\left(\frac{-1}{r-1}\right)^{|T|}
$$

Then (7) now gives that

$$
E(X)=\frac{-\alpha^{2}}{\alpha-\alpha^{2}}=\frac{-\alpha}{1-\alpha} .
$$

We now break our analysis into three cases

- $\alpha>1 / r$. We will see that this cannot happen, which implies Claim 4.1.
- $\alpha=1 / r$. The analysis of this case will imply Theorem 1.1.
- $\alpha=1 / r-\varepsilon$. This case leads to the proof of Theorem 1.2.

For all three cases it is important to note that for all $T, X(T)=X \geq \frac{-1}{r-1}$ with equality iff $|T|=1$.

1. Let $\alpha>\frac{1}{r}$. Then $E(X)=\frac{-\alpha}{1-\alpha}<\frac{-1}{r-1}$. However this is impossible, because $X \geq \frac{-1}{r-1}$. Hence if $I$ is an independent set then $\mu(I) \leq \frac{1}{r}$. Thus we have just proven Claim 4.1 (again).
2. Next, if $\alpha=1 / r$ then $E(X)=\frac{-1}{r-1}$ which implies that $X \equiv \frac{-1}{r-1}$ meaning that all the support of $|\widehat{f}(S)|^{2}$ must be concentrated on $u_{S}$ with $|S|=1$. Hence $f$ has all its Fourier transform concentrated on the first two levels. According to Lemma 2.3 this implies that $f$ depends on one coordinate only, and hence, since $f$ is not constant, it is the indicator function of $S_{l}=k$ for some $l$ and $k$, as asserted by Theorem 1.1.
3. Finally we analyze the case where $\alpha$ is slightly less than $1 / r$. This means that $E(X)$ is very close to $\frac{-1}{r-1}$, its minimal value. Observing that for $S,|S|>1$,

$$
X(S) \geq \frac{-1}{(r-1)^{3}}>\frac{-1}{r-1}
$$

this implies that almost all the weight of $|\widehat{f}(S)|^{2}$ is concentrated on $S$, $|S| \leq 1$.
Formally, let

$$
Y=X+\frac{1}{r-1} \geq 0
$$

When $Y>0$ then $Y \geq \frac{-1}{(r-1)^{3}}+\frac{1}{r-1}=\frac{r(r-2)}{(r-1)^{3}}$. Hence, by Markov's inequality,

$$
\begin{equation*}
\operatorname{Prob}[Y>0] \leq E(Y) \cdot \frac{(r-1)^{3}}{r(r-2)} \tag{8}
\end{equation*}
$$

Now, if $\alpha>1 / r-\varepsilon$ then

$$
E(Y)=E(X)+\frac{1}{r-1}<\frac{\varepsilon r^{2}}{(r-1)(r \varepsilon+r-1)}
$$

which, using (8) together with $r \geq 3$, implies

$$
\begin{equation*}
\operatorname{Prob}[Y>0] \leq \frac{\varepsilon r}{r \varepsilon+r-1} \cdot \frac{(r-1)^{2}}{(r-2)} \leq 2 r \varepsilon \tag{9}
\end{equation*}
$$

Recalling that

$$
Y(S)>0 \quad \text { iff } \quad X(S)>\frac{-1}{r-1} \text { iff }|S|>1
$$

the upshot of (9) is that $\forall \varepsilon>0$ and for every independent set $I$ with $\mu(I)>1 / r-\varepsilon$ we have

$$
\sum_{|S|>1}|\widehat{f}(S)|^{2}=\left(\alpha-\alpha^{2}\right) \operatorname{Prob}[Y>0] \leq 2 \varepsilon
$$

where $f$ is the characteristic function of $I$. We now deduce directly from Lemma 2.4 that there is some 'dictator' function $g$ that depends on at most one coordinate, and such that $\|f-g\|_{2}^{2}<\frac{K}{\alpha-\alpha^{2}-\varepsilon} \cdot 2 \varepsilon=O(\varepsilon)$. (Here, if $\varepsilon$ is not sufficiently small, $g \equiv 0$ is also a possibility, but if $\varepsilon$ is such that $\|f\|_{2}^{2}-\frac{K}{\alpha-\alpha^{2}-\varepsilon} \cdot 2 \varepsilon>0$ then $g$ must be a dictator.) $\quad \square$
4.1 Infinite weak products. Let $A$ be an infinite set and consider $G=\times{ }^{A} K_{r}$, the infinite weak product of $K_{r}$. Let $\mu$ be the product measure on $A$, the measure generated by the cylinders of finite codimension. The set $\left\{u_{S}\right\}_{|S|<\infty}$ is a basis for the measurable functions, hence the same proof technique as before yields the following:
Theorem 4.4. Let $I \subseteq \times{ }^{A} K_{r}$ be a measurable independent set with $\mu(I)=1 / r$. Then there exists $j \in A$ and a set $J \subset G$ of the form

$$
J=\left\{S: S_{j}=k\right\}
$$

with

$$
\mu(J \triangle I)=0
$$

Remarks. - Alternatively one can deduce the infinite dimensional statement from the finite case using the stability statement of Theorem 1.2. It seems impossible to make this deduction directly without the stability theorem.

- In [GL], an example due to Pósa is presented that shows, using ultrafilters, an optimal coloring of $G$ which is not of the form $\sigma(S)=S_{j}$. Of course, the color classes thus formed must be non-measurable.
4.2 Large independent sets in weak products. Theorem 1.1 gives a complete characterization of independent sets of measure $1 / r$. Theorem 1.2 also characterizes independent sets of size close to $1 / r$. A natural question is whether smaller independent sets, of size, say, $1 /(100 r)$ still have some simple structure. It seems quite hopeless to approach this question using traditional combinatorial arguments. We now present a conjecture regarding such sets, that may be susceptible to Fourier methods.

Let us start with an example. Consider $G=\times{ }^{n} K_{3}$ over vertices $\{0,1,2\}^{n}$. Let
$I=\left\{S\right.$ : at least two of the three coordinates $S_{1}, S_{2}, S_{3}$ are equal to 0$\}$.
Clearly, if $S, T \in I$ then they must agree on one of the first three coordinates, hence $I$ is an independent set, with $\mu(I)=7 / 27$. A naive conjecture would be that any reasonably large independent set is similarly determined by few coordinates. However, this is clearly too strong an assumption, since a random subset of $I$ is also an independent set. Moreover, this problem does not go away if we consider only maximal independent sets (with respect to containment). Indeed consider any 'nice' independent set $I_{0}$, e.g. the one above with the following incremental modification. At each step select some random vertex $v \in\{0,1,2\}^{n}$ and let $I_{t}$ be an arbitrarily chosen maximal independent set containing $\left(I_{t-1} \cup\{v\}\right) \backslash \Gamma(v)$ where $\Gamma(v)$ are the neighbors of $v$ in $G$. By definition, at each step $I_{t}$ remains a maximal independent set, and after sufficiently many steps its size is still about half that of $I$, and it certainly is not determined by only a few coordinates.

Thus, we conjecture that every large independent set is close to being contained in a set determined by few coordinates. Sets of the form $\left\{S: S_{j}=k\right\}$ are sometimes referred to as dictatorships (dictated by the $j$ th coordinate.) Analogously, we call $S$ a $k$-junta if its characteristic function depends on at most $k$ different coordinates.

Conjecture 4.5. There exists a function $M=M(r, \alpha, \varepsilon)$ (note that $M$ does not depend on $n!$ ), such that for every integer $r \geq 3$ and for every $\alpha>0$ if $I \subseteq K_{r}^{n}$ is an independent set with $\mu(I)>\alpha$ then for any $\varepsilon>0$ there is an $M$-junta $J$ which is itself an independent set such that

$$
\mu(I \backslash J)<\varepsilon,
$$

i.e. I is almost contained in a junta-type independent set.

We believe that this conjecture, if proven, may be quite useful in resolving the hardness of approximate-coloring 3 -colorable graphs, an infamously open hardness of approximation problem. We outline the connection roughly as follows. Following the PCP approach we begin with some NP-hard constraint satisfaction problem (CSP) and construct from a given CSP instance a 3 -colorable graph. This graph would consist of many copies of $\left(K_{3}\right)^{n}$ that are further interconnected, such that each copy represents a variable in the CSP taking values in $\{1, \ldots, n\}$. One would use the above conjecture to prove that a given coloring of this graph corresponds, in each copy of $\left(K_{3}\right)^{n}$, to a 'commitment' assigning some constant number of values in $[n]$ to that variable (these are the coordinates of the junta). Then, the structure of the graph would be used to show that these values must satisfy the constraints of the CSP we started out with.

This scheme has been successfully used with similar approximation problems such as vertex-cover [DS], where similarly structured graphs containing many interconnected copies of $\{0,1\}^{n}$ have been constructed, and combinatorial statements regarding the 'juntativity' of subsets of $\{0,1\}^{n}$ were essential. One must add that it is not altogether clear how to construct such a reduction and such a graph in the 3 -colorability case, so proving the above conjecture is by no means the last step of completing such a proof.

Just as the generalization of the results from [FKN], which characterize Boolean functions that are close to dictatorships, are essential for proving Theorem 1.2, we believe that for the proof of the above conjecture it will be necessary to generalize the junta-type results of [Bou2] and [KiS]. First steps towards this goal appear in [DFK].

## 5 Independent Sets in $\boldsymbol{H}^{\boldsymbol{n}}$

We begin this section with the simple proof of Proposition 1.3. Let $H$ be a vertex-transitive graph on a set $V=\left\{v_{1}, \ldots, v_{r}\right\}$ of $r$ vertices, let $A$ be the automorphism group of $H$, and let $S$ be an independent set of
maximum size in $G=H^{n}$. Let $a_{1}, a_{2}, \ldots, a_{n}$ be $n$ random elements of $A$, where each $a_{i}$ is chosen, randomly, independently and uniformly, among the members of $A$. For each $i, 1 \leq i \leq r$, let $x_{i}$ be the random vertex of $G$ given by $x_{i}=\left(a_{1}\left(v_{i}\right), a_{2}\left(v_{i}\right), \ldots, a_{n}\left(v_{i}\right)\right)$. It is easy to see that $x_{i}$ is uniformly distributed among all vertices of $G$. Therefore, the expected size of the intersection $S \cap\left\{x_{1}, \ldots, x_{r}\right\}$ is precisely $|S| \frac{r}{|G|}$. Note, however, that the induced subgraph of $G$ on $\left\{x_{1}, \ldots, x_{r}\right\}$ is a copy of $H$, and therefore the size of the intersection is at most $\alpha(H)$ for every choice of $a_{1}, a_{2}, \ldots, a_{n}$. We conclude that $|S| \frac{r}{|G|} \leq \alpha(H)$, and as $|S|=\alpha(G)$ this implies that

$$
\frac{\alpha(G)}{|G|} \leq \frac{\alpha(H)}{|H|}
$$

as needed.
REmARK. The transitivity assumption is needed in the above proposition, that is, its assertion fails for some connected regular base graphs $H$. Indeed, consider, for example, the following $d$-regular graph $H$. Its vertex set is the disjoint union of three sets, $A, B$ and $C$, where $|A|=|B|=d$ and $|C|=$ $d+1$. Let $a, b, c_{1}, c_{2}$ be distinguished elements of the corresponding sets, and let the edge set of $H$ consist of all edges between $A$ and $B$ except $a b$, all internal edges in $C$ except $c_{1} c_{2}$, and the two additional edges $a c_{1}$ and $b c_{2}$. It is easy to check that $\alpha(H)=d+1$. For each $n \geq 1$, let $m_{n}=\alpha\left(H^{n}\right) /\left|H^{n}\right|$ denote the maximum possible measure of an independent set in $H^{n}$. Thus $m_{1}=\frac{d+1}{3 d+1}$. We claim that for $n>1$,

$$
m_{n} \geq \frac{d}{3 d+1}+\frac{d}{3 d+1} m_{n-1}
$$

To prove this claim, note that the set of vertices of $H^{n}$ whose first coordinate is a member of $A$ (with all possible extensions to the rest of the coordinates), together with the set of all vertices of $H^{n}$ whose first coordinate is in $C-\left\{c_{1}\right\}$, with the last $n-1$ coordinates belonging to a maximum size independent set in $H^{n-1}$, form an independent set in $H^{n}$. This proves the claim and implies that

$$
m_{n} \geq \frac{d}{3 d+1}+\left(\frac{d}{3 d+1}\right)^{2}+\cdots+\left(\frac{d}{3 d+1}\right)^{n-1}+\left(\frac{d}{3 d+1}\right)^{n-1} \frac{d+1}{3 d+1}
$$

For any fixed $d$, the above lower bound tends to $\frac{d}{2 d+1}$ as $n$ tends to infinity, showing that the graph $H$ does not satisfy the assertion of Proposition 1.3.

We proceed with the proofs of Theorems 1.4 and 1.5. These are based on some basic facts about graph eigenvalues and eigenvectors. The relation
between the spectral properties of a graph and its structural properties has been studied extensively, cf., e.g. [A1], $[\mathrm{KrS}]$ and its references.

The adjacency matrix of a graph $H=(V, E)$ is the matrix $A=\left(a_{u, v}\right)_{u, v \in V}$, in which $a_{u, v}=1$ if $u v \in E$ and $a_{u, v}=0$ otherwise. This matrix is symmetric and hence has real eigenvalues $\mu_{1} \geq \ldots \geq \mu_{r}$ (called the eigenvalues of $H$ ) and an orthonormal basis of eigenvectors $u_{1}, \ldots, u_{r}$. If $H$ is connected and $d$-regular, then $\mu_{1}=d$, and the corresponding (normalized) eigenvector is $u_{1}=\frac{1}{\sqrt{r}}(1,1, \ldots, 1)^{t}$. By the Perron-Frobenius theorem $d \geq\left|\mu_{i}\right|$ for all $i$, and $\mu_{2}<d$ iff $H$ is connected, whereas $\mu_{r}>-d$ iff $H$ is non-bipartite. Note that the trace of $A$ is 0 , as $H$ has no loops, and therefore if $H$ has edges, then $\mu_{1}$ is strictly positive and $\mu_{r}$ is strictly negative.

The adjacency matrix of the weak power $H^{n}$ of $H$ is the tensor $n$-th power of $A$. Its eigenvalues are all products $\prod_{i=1}^{n} \mu_{g(i)}$, as $g$ ranges over all $r^{n}$ functions mapping $\{1,2, \ldots, n\}$ into $\{1,2, \ldots, r\}$. The corresponding orthonormal basis of eigenvectors is obtained by taking all tensor products of the eigenvectors $u_{i}$.
Lemma 5.1. Let $G=(V, E)$ be a $D$-regular, connected graph on $m$ vertices, let $D=\lambda_{1} \geq \lambda_{2} \geq \ldots \geq \lambda_{m}$ be its eigenvalues, and let $v_{1}, v_{2}, \ldots, v_{m}$ be the corresponding orthonormal basis of eigenvectors, where $v_{1}=\frac{1}{\sqrt{m}}(1,1, \ldots, 1)^{t}$. Let $I$ be an independent set in $G$, let $f$ be its characteristic function, and suppose $f=\sum_{i=1}^{m} c_{i} v_{i}$.
(i) The size of I satisfies

$$
\begin{equation*}
|I| \leq \frac{-\lambda_{m} m}{D-\lambda_{m}} \tag{10}
\end{equation*}
$$

and equality holds if and only if $f$ is a linear combination of the eigenvectors corresponding to the eigenvalues $\lambda_{1}$ and $\lambda_{m}$.
(ii) If $\sum_{i: \lambda_{i} \neq \lambda_{1}, \lambda_{m}} c_{i}^{2}=\delta m$ and $\lambda=\min \left\{\lambda_{i}: \lambda_{i} \neq \lambda_{m}\right\}$ then

$$
\begin{equation*}
\delta \leq \frac{\left(D-\lambda_{m}\right)}{\left(\lambda-\lambda_{m}\right)} \frac{|I|}{m}\left(\frac{-\lambda_{m}}{D-\lambda_{m}}-\frac{|I|}{m}\right) \tag{11}
\end{equation*}
$$

Proof. Let $A$ be the adjacency matrix of $G$. Since $f$ is the characteristic vector of an independent set, $f^{t} A f=0$. As the vectors $v_{i}$ form an orthonormal basis, $f^{t} A f=\sum_{i=1}^{m} c_{i}^{2} \lambda_{i}, \sum_{i=1}^{m} c_{i}^{2}=f^{t} f=\sum_{v \in V} f^{2}(v)=|I|$, and $c_{1}=f^{t} v_{1}=|I| / \sqrt{m}$. It thus follows that

$$
\begin{aligned}
0=f^{t} A f=\sum_{i=1}^{m} c_{i}^{2} \lambda_{i} & =D c_{1}^{2}+\sum_{i=2}^{m} c_{i}^{2} \lambda_{i} \geq D \frac{|I|^{2}}{m}+\lambda_{m} \sum_{i=2}^{m} c_{i}^{2} \\
& =D \frac{|I|^{2}}{m}+\lambda_{m}\left(|I|-\frac{|I|^{2}}{m}\right)
\end{aligned}
$$

Therefore,

$$
-\lambda_{m}|I| \geq \frac{|I|^{2}}{m}\left(D-\lambda_{m}\right)
$$

implying (10). Moreover, equality holds if and only if $c_{i}=0$ whenever $\lambda_{i}$ is neither $D=\lambda_{1}$ nor $\lambda_{m}$. This establishes the assertion of part (i).

The proof of part (ii) is very similar. Here

$$
\begin{aligned}
0=f^{t} A f=\sum_{i=1}^{m} c_{i}^{2} \lambda_{i} & =D \frac{|I|^{2}}{m}+\lambda_{m}\left(|I|-\frac{|I|^{2}}{m}-\delta m\right)+\sum_{i: \lambda_{i} \neq \lambda_{1}, \lambda_{m}} c_{i}^{2} \lambda_{i} \\
& \geq D \frac{|I|^{2}}{m}+\lambda_{m}\left(|I|-\frac{|I|^{2}}{m}-\delta m\right)+\delta m \lambda .
\end{aligned}
$$

Therefore,

$$
\left(D-\lambda_{m}\right)|I|\left(\frac{-\lambda_{m}}{D-\lambda_{m}}-\frac{|I|}{m}\right)=-\lambda_{m}|I|-\frac{D-\lambda_{m}}{m}|I|^{2} \geq \delta m\left(\lambda-\lambda_{m}\right)
$$

implying (11).
Proof of Theorem 1.4. Put $G=H^{n}$ and let $m=r^{n}$ be the number of vertices of $G$. By the discussion preceding Lemma 5.1, the eigenvalues of $G$ are all products of eigenvalues of $H$ and the corresponding eigenvectors are the tensor products of the eigenvectors of $H$. In particular, the largest eigenvalue of $G$ is its degree $\lambda_{1}=D=d^{n}$, and its smallest (most negative) eigenvalue is $\lambda_{m}=\mu_{r} d^{n-1}$. Therefore, by Lemma 5.1, part (i),

$$
\alpha(G) \leq \frac{-\lambda_{m} m}{D-\lambda_{m}}=\frac{-\mu_{r} d^{n-1} r^{n}}{d^{n}-\mu_{r} d^{n-1}}=\frac{-\mu_{r} r^{n}}{d-\mu_{r}}
$$

as needed. Moreover, if $I$ is an independent set of size $\frac{-\mu_{r} r^{n}}{d-\mu_{r}}$ in $G$ then, by Lemma 5.1 part (i), the characteristic function $f$ of $I$ is a linear combination of the eigenvectors corresponding to the eigenvalues $\lambda_{1}$ and $\lambda_{m}$. Since $H$ is connected and non-bipartite the multiplicity of the eigenvalue $\lambda_{1}=d^{n}$ of $G$ is one, and its only eigenvector is the vector $\frac{1}{r^{n / 2}}(1,1, \ldots, 1)^{t}$ obtained by taking the tensor $n$-th power of the first eigenvector of $H$. Similarly, the only eigenvectors corresponding to the smallest eigenvalue $\lambda_{m}=\mu_{r} d^{n-1}$ are the tensor products of one eigenvector of $H$ corresponding to its smallest eigenvalue, with $n-1$ copies of the vector $\frac{1}{\sqrt{r}}(1,1, \ldots, 1)^{t}$. Let $u_{1}, u_{2}, \ldots, u_{r}$ be an orthonormal basis of eigenvectors of $H$, where $u_{i}$ corresponds to the eigenvalue $\mu_{i}$. In particular, $u_{1}=\frac{1}{\sqrt{r}}(1,1, \ldots, 1)^{t}$. For each sequence $S \in\{1,2, \ldots, r\}^{n}$, let $v_{S}$ denote the tensor product $\prod_{i=1}^{n} u_{S_{i}}$. Define, also $|S|=\left|\left\{i: S_{i} \neq 1\right\}\right|$. In addition, if $|S|=1$, let $i(S)$ denote the unique index $i$ such that $S_{i} \neq 1$. The vectors $v_{S}$ form an orthonormal basis of eigenvectors of $G$, and by the above discussion the vector $f$ is a linear combination of vectors $v_{S}$ with $|S| \leq 1$. Let $f=\sum_{|S| \leq 1} \alpha_{S} v_{S}$ be this combination.

Imitating the proof of Lemma 2.3 we next show that since $f$ is Boolean, all the vectors $S$ for which $|S|=1$ and $\alpha_{S} \neq 0$ have the same value of $i(S)$. Indeed, since $f$ is Boolean, $f^{2}=f$. However, $f^{2}=\sum_{S, S^{\prime}} \alpha_{S} \alpha_{S^{\prime}} v_{S} v_{S^{\prime}}$, where here $v_{S} v_{S}^{\prime}$ denotes the vector whose coordinates are the products of the corresponding coordinates of $v_{S}$ and $v_{S^{\prime}}$. It is easy to check that if $|S|=\left|S^{\prime}\right|=1$ and $i(S)=i\left(S^{\prime}\right)$, then $v_{S} v_{S^{\prime}}$ is a linear combination of the vectors $v_{T}$ with $i(T)=i(S)$ and the constant vector. It is also easy to check that for the unique sequence $S=\overline{1}=(1,1, \ldots, 1)$, and for any $S^{\prime}$ with $\left|S^{\prime}\right|=1, v_{\overline{1}} v_{S^{\prime}}$ is a scalar multiple of $v_{S^{\prime}}$. On the other hand, if $|S|=\left|S^{\prime}\right|=1$ and $i(S) \neq i\left(S^{\prime}\right)$, then $v_{S} v_{S^{\prime}}$ is a multiple of the vector $v_{T}$, where $T$ coincides with $S$ in coordinate number $i(S)$, with $S^{\prime}$ in coordinate number $i\left(S^{\prime}\right)$, and with both in all other coordinates (in which they are both 1 ). In particular, $|T|=2$. By the unique representation of $f$ as a linear combination of the members of the basis $v_{S}$, it thus follows that all vectors $S$ for which $\alpha_{S} \neq 0$ and $|S|=1$ have the same value of $i(S)$. In other words, $f$ is determined by a single coordinate, that is, there is a coordinate $i$ such that for each vertex $x=\left(x_{1}, x_{2}, \ldots, x_{n}\right)$ of $H^{n}, x$ lies in $I$ depending only on the value of $x_{i}$. It follows that all the vertices that appear in the $i$-th coordinate of members of $I$ form an independent set of $H$, and as this set has to be of maximum size, the assertion of the second part of Theorem 1.4 follows.

Proof of Theorem 1.5. Let $f$ be the characteristic vector of $I$, put $D=d^{n}, m=r^{n}$ and let $\lambda_{1}=D \geq \lambda_{2} \geq \cdots \geq \lambda_{m}$ be the eigenvalues of $G$. Define also $\lambda=\min \left\{\lambda_{i}: \lambda_{i} \neq \lambda_{m}\right\}$. Let $u_{1}, \ldots, u_{r}$ as well as $\left\{v_{S}: S \in\{1,2, \ldots, r\}^{n}\right\}$ be, as in the previous proof, the eigenvectors of $H$ and of $G$, respectively. Let $f=\sum c_{S} v_{S}$ be a representation of $f$ as a linear combination of the vectors $v_{S}$, and define $\delta$ by

$$
\delta m=\sum_{|S|>1} c_{S}^{2}
$$

By Lemma 5.1, part (ii), we conclude that

$$
\begin{equation*}
\delta \leq \frac{\left(D-\lambda_{m}\right)}{\left(\lambda-\lambda_{m}\right)} \frac{|I|}{m} \epsilon \tag{12}
\end{equation*}
$$

It is not difficult to express $\lambda$ as the minimum between three possible eigenvalues of $G$. If $\mu=\min \left\{\mu_{i}: \mu_{i}>\mu_{r}\right\}$, then $\lambda=\min \left\{\mu d^{n-1}, \mu_{r} \mu_{2} d^{n-2}, \mu_{r}^{3} d^{n-3}\right\}$. In any case, since $\mu_{2}<d$ and $\mu_{r}>-d$, there exists some $\eta=\eta(H)>0$ (which depends only on the base graph $H$, and not on $n$ ), such that $\lambda=\left(\mu_{r}+\eta\right) d^{n-1}$. Substituting in (12) and recalling that $D=d^{n}$ and
$\lambda_{m}=\mu_{r} d^{n-1}$ we conclude that

$$
\delta \leq \frac{\left(d-\mu_{r}\right)}{\eta} \frac{|I|}{m} \epsilon .
$$

Now, for any function $g$ let $(\widehat{g})$ be the vector of coefficients of its expansion according to the characters $u_{S}$ described in section 2 , and ( $\bar{g}$ ) the vector of coefficients of its expansion according to the set $v_{S}$. Since each of these two sets is, up to a multiplicative constant, an orthonormal basis the transformation between these two expansions is norm-preserving up to a multiplicative constant:

$$
m\|(\widehat{g})\|_{2}^{2}=\|(\bar{g})\|_{2}^{2}
$$

Furthermore, what is crucial to us, is that this transformation is normpreserving levelwise: for any $0 \leq k \leq n$ the $l_{2}$ norm of $\widehat{g}$ on the $k$ 'th level is precisely $\sqrt{m}$ times the $l_{2}$ norm of $\bar{g}$ on the $k^{\prime}$ 'th level:

$$
m \sum_{|S|=k}|\widehat{g}(S)|^{2}=\sum_{|S|=k}|\bar{g}(S)|^{2} .
$$

This is true because for any such $k$

$$
\operatorname{Span}\left(\left\{u_{S}:|S|=k\right\}\right)=\operatorname{Span}\left(\left\{v_{S}:|S|=k\right\}\right),
$$

giving a simultaneous level-wise decomposition of the space of functions into orthogonal subspaces.

Returning to our function $f$ we perform this change of basis, and express $f$ as a linear combination of the characters $u_{S}$ and conclude that

$$
\sum_{|S|>1}|\widehat{f}(S)|^{2}=\frac{\sum_{|S|>1} c_{S}^{2}}{m}=\frac{\delta m}{m} \leq \frac{\left(d-\mu_{r}\right)}{\eta} \frac{|I|}{m} \epsilon .
$$

The desired result now follows from Lemma 2.4.
5.1 Examples. There are various examples of graphs $H$ which satisfy the assumptions of Theorems 1.4 and 1.5. These include complete graphs, line graphs of regular graphs which contain a perfect matching, Kneser graphs, some strongly regular graphs and appropriate classes of random regular graphs. Here are the details:

- Let $k \geq 3$ and let $G=(V, E)$ be a $k$-regular graph on $n$ vertices which contains a perfect matching, i.e. a set of $n / 2$ pairwise disjoint edges. The line graph $L(G)$ of $G$ has a vertex for every edge of $G$ and two vertices of $L(G)$ are adjacent if as edges of $G$ they share a vertex. By definition, $L(G)$ is a $d=2(k-1)$-regular graph on $r=n k / 2$ vertices. Since an independent set in $L(G)$ corresponds to a matching in $G$
we have that $\alpha(L(G))=n / 2$. As noted by A. Hoffman, the smallest eigenvalue of the adjacency matrix $A_{L(G)}$ of $L(G)$ is easy to compute. Let $B=\left(b_{v, e}\right)$ be the vertex-edge incidence matrix which is defined as $b_{v, e}=1$ iff $v \in e$. Then $A_{L(G)}=B^{T} B-2 I$. Since $B^{T} B$ is positive semidefinite, all its eigenvalues are non-negative. Hence, the eigenvalues of $A_{L(G)}$ are $\geq-2$. Moreover, we have that the rank of $B^{T} B$ satisfies $r\left(B^{T} B\right)=r(B) \leq|V(G)|<|E(G)|$. So, $B^{T} B$ has at least one 0 eigenvalue. This implies that the smallest eigenvalue of $A_{L(G)}$ is $\mu_{r}=-2$. Now it is easy to check that $L(G)$ satisfies the assumptions of Theorems 1.4 and 1.5. If, in addition, $G$ is class 1 -Vizing (that is, it is $k$-edge colorable), then $H=L(G)$ satisfies the assumption in the remark following Theorem 1.4, and hence any optimal coloring of $H^{n}$ is determined by one of the factors.
- Let $t, k$ be two positive integers, $t>2 k$. The Kneser graph $K(t, k)$ is the graph whose vertices are all $k$-subsets of a $t$-element set, where two vertices are adjacent iff the corresponding sets are disjoint. The number of vertices of this graph is clearly $r=\binom{t}{k}$, and it is well known that its eigenvalues are all numbers of the form $(-1)^{j}\binom{t-k-j}{k-j}$, for $j \in\{0,1, \ldots, k\}$ (see, e.g. [Lo] for a proof.) In particular, the largest eigenvalue is the degree of regularity $d=\mu_{1}=\binom{t-k}{k}$ and the smallest is $\mu_{r}=-\binom{t-k-1}{k-1}$. The independence number of $K(t, k)$ is the maximum possible cardinality of an intersecting family of $k$-subsets of a $t$-element set. This maximum is determined by the Erdős-Ko-Rado theorem [EKR], and is the number of all $k$-sets containing a single element, that is

$$
\binom{t-1}{k-1}=\frac{\binom{t-k-1}{k-1}\binom{t}{k}}{\binom{t-k}{k}+\binom{t-k-1}{k-1}}=\frac{-\mu_{r} r}{d-\mu_{r}} .
$$

Therefore, for each admissible $t$ and $k, H=K(t, k)$ satisfies the assumptions of Theorems 1.4 and 1.5. In particular, $K(5,2)$, which is the Petersen graph, satisfies these assumptions.

- Let $a$ be an even integer and let $q=p^{a}$ be a prime power which is congruent to 1 modulo 4 so that -1 is a square in the finite field $G F(q)$. Let $P_{q}$ be the graph whose vertices are all elements of $G F(q)$ and two vertices are adjacent if and only if their difference is a quadratic residue in $G F(q)$. This graph is usually called the Paley graph. It is easy to see that $P_{q}$ is $d=(q-1) / 2$-regular. To compute the eigenvalues of $P_{q}$ one can use the fact that it is strongly regular, i.e. the number of common neighbors of any two vertices of it has the same
value - $(q-5) / 4$ if they are adjacent, and the same value - $(q-1) / 4$ if they are not. This implies that the smallest eigenvalue of the adjacency matrix of $P_{q}$ is $\mu_{q}=-(\sqrt{q}+1) / 2$ (see, e.g. $\left.[\mathrm{KrS}]\right)$. It is easy to see that all elements of the subfield $G F(\sqrt{q}) \subset G F(q)$ are quadratic residues in $G F(q)$. This implies that for every quadratic non-residue $\beta \in G F(q)$ all elements of any multiplicative coset $\beta G F(\sqrt{q})$ form an independent set of size $\sqrt{q}$. Hence for this graph we get

$$
\alpha\left(P_{q}\right)=\sqrt{q}=\frac{(\sqrt{q}+1) / 2}{(q-1) / 2+(\sqrt{q}+1) / 2} q=\frac{-\mu_{q}}{d-\mu_{q}} q .
$$

The graph $P_{q}$ also satisfies the condition in the remark following Theorem 1.4. Indeed, it is easy to see that the sets $\{\alpha+\beta G F(\sqrt{q})$ : $\alpha \in G F(\sqrt{q})\}$ form a partition of of its set of vertices into $\sqrt{q}$ independent sets.

- Another construction of a strongly regular graph which satisfies the assumptions of Theorems 1.4 and 1.5 is due to Delsarte and Goethals and to Turyn (see, e.g. $[\mathrm{KrS}]$ ). Let $q$ be a prime power and let $V(H)$ be the elements of the two dimensional vector space over $G F(q)$. Thus $H$ has $r=q^{2}$ vertices. Partition the $q+1$ lines through the origin of the space into two sets $P$ and $N$, where $|P|=k$. Two vertices $x$ and $y$ of the graph $H$ are adjacent if $x-y$ is parallel to a line in $P$. It is easy to check that $H$ is $d=k(q-1)$-regular. Since all the points on any line from $P$ form an independent set in $H$, we have $\alpha(H) \geq q$ and $\chi(H) \leq q$. It is known that the smallest eigenvalue of $H$ is $\mu_{r}=-k$. Therefore, by Lemma 5.1, we obtain that

$$
q \leq \alpha(H) \leq \frac{-\mu_{r}}{d-\mu_{r}} r=\frac{k}{k(q-1)+k} q^{2}=q,
$$

showing that $H$ satisfies the assumptions of Theorems 1.4 and 1.5, as well as that of the remark between them.

- Let $k \geq 3$ be a fixed integer, let $r$ be an integer divisible by $k$, and let $d<r(k-1) / k$ be another integer, divisible by $k-1$. Let $H$ be a random $k$-colorable graph obtained as follows. The vertex set of $H$ is the union of $k$ pairwise disjoint sets $A_{1}, A_{2}, \ldots, A_{k}$, each of size $r / k$. The edges of $H$ are the union of $\frac{d}{k-1}\binom{k}{2}$ random matchings $M_{i j s}$, where for $1 \leq i<j \leq k$, and $1 \leq s \leq d /(k-1), M_{i j s}$ is a random perfect matching between $A_{i}$ and $A_{j}$. The resulting graph $H$ is obviously $d$ regular, and it is not difficult to see that $-d /(k-1)$ is an eigenvalue of it with multiplicity at least $k-1$. In fact, each vector which is constant on the members of each $A_{i}$, where the sum of these $k$ con-
stants is 0 , is an eigenvector of this eigenvalue. As shown in [AK], for any fixed $k$, and all $d>d_{0}(k)$ and $r>r_{0}(k, d)$, all the other eigenvalues of $H$ besides $d$ and $-d /(k-1)$ are in absolute value at most $O(\sqrt{d})$ with high probability. It follows that with high probability $H$ satisfies the assumptions of Theorems 1.4 and 1.5 , as well as the assumption in the remark between them. Another, similar class of graphs satisfying these assumptions can be constructed by replacing each bipartite graph between $A_{i}$ and $A_{j}$ in the construction above by an appropriate $d /(k-1)$-regular pseudo-random graph. We omit the detailed description.


## 6 Concluding Remarks

- The weak product considered here is similar to the AND product considered in the definition of the Shannon capacity of graphs. The $n$-th $A N D$ power of an undirected graph $H=(V, E)$ is the graph denoted by $H^{\wedge n}$ whose vertex set is $V^{n}$ in which distinct vertices $\left(x_{1} \ldots x_{n}\right)$ and $\left(x_{1}^{\prime} \ldots x_{n}^{\prime}\right)$ are connected if $\left\{x_{i}, x_{i}^{\prime}\right\} \in E$ for all $i \in\{1,2, \ldots, n\}$ such that $x_{i} \neq x_{i}^{\prime}$. The Shannon capacity of $H$ is the limit of $\left(\alpha\left(H^{\wedge n}\right)\right)^{1 / n}$, as $n$ tends to infinity. Although the definition here is very similar to that of a weak product, the investigation of the independence numbers of these powers is far more complicated than that of the independence numbers in weak powers. In fact, even the Shannon capacity of a cycle of length 7 is not known. See [Lo], [A2] and their references for more details.
- The $n$-th (sparse) power of an undirected graph $H=(V, E)$, denoted here by $H^{* n}$, is the graph whose vertex set is $V^{n}$ in which distinct vertices $\left(x_{1} \ldots x_{n}\right)$ and $\left(x_{1}^{\prime} \ldots x_{n}^{\prime}\right)$ are connected iff there exists a single index $i$ such that $x_{j}=x_{j}^{\prime}$ for all $j \neq i$ and $x_{i}$ and $x_{i}^{\prime}$ are connected in $H$. There are several papers dealing with the asymptotic behaviour of the independence number of $H^{* n}$, for a fixed graph $H$. See, e.g. [HHP] and its references. In particular, it is known that for every fixed graph $H$, the limit $\lim _{n \mapsto \infty} \alpha\left(H^{* n}\right) /|V|^{n}$ exists, and is at least $1 / \chi(H)$ and at most the reciprocal of the fractional chromatic number of $H$. Moreover, for every Cayley graph of an Abelian group this limit is precisely the reciprocal of the fractional chromatic number of $G$.
- Most of the results proved here can be extended to independent sets in weak products of different regular factors $H_{i}$, and there is no need to assume all the factors are identical. We omit the detailed statements of the corresponding results.
- The assumptions in Theorems 1.4 and 1.5 that the base graph $H$ is connected and non-bipartite are essential. Indeed, if $H$ is either disconnected or bipartite, then $H^{n}$ has many connected components, and as in each of them one can choose a maximum independent set separately, no single coordinate can determine the membership in a maximum independent set of $H^{n}$. Moreover, if $H$ is bipartite then so is $H^{n}$, and as it is regular, its independence number is precisely half the number of vertices. There are many maximum size independent sets. In fact, if $J$ is a maximum size independent set in $H,|J|=|H| / 2$, and $A \subset\{1,2, \ldots, n\}$ is any set of odd size, then the set

$$
I=\left\{\left(x_{1}, \ldots, x_{n}\right): \mid\left\{i: i \in A \text { and } x_{i} \in J\right\}|>|A| / 2\}\right.
$$

is a maximum size independent set in $H^{n}$. Therefore, there are no uniqueness or stability results in this case. In case $H$ is disconnected, one can consider each component separately, and use the assertions of Theorems 1.4 and 1.5 for products of possibly distinct factors to obtain some information about the structure of maximum and near maximum independent sets in $H^{n}$.

- The proof of Theorem 1.5 works even if the set $I$ considered is not independent, but contains a small number of edges. Indeed, the assertion of Lemma 5.1, part (ii), holds, with an extra constant factor, even if the induced subgraph of $G$ on $I$ contains some $\delta m\left(\lambda-\lambda_{m}\right) / 4$ edges. This can be used to obtain results about colorings of $H^{n}$ in which the number of monochromatic edges is relatively small.
- The proof of Theorem 1.5 is based on the fact that a change of basis enables us to prove several results from discrete Harmonic Analysis for general tensor products of an arbitrary fixed orthogonal basis, and not only for characters. We believe that this simple fact is likely to lead to additional interesting consequences.

Note added in proof. In a recent paper [LT], Larose and Tardif study the structure of maximum independent sets in graph powers. They prove Proposition 1.3, and describe an example of vertex transitive graphs $H$ for which the assertion of Theorem 1.4 fails.

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Noga Alon, Schools of Mathematics and Computer Science, Sackler Faculty of Exact Sciences, Tel Aviv University, Tel Aviv, Israel noga@math.tau.ac.il

Irit Dinur, NEC Research, Princeton, NJ 08540, USA iritd@nec-labs.com
Ehud Friedgut, Institute of Mathematics, Hebrew University, Jerusalem, Israel ehudf@math.huji.ac.il
Benny Sudakov, Department of Mathematics, Princeton University, Princeton, NJ 08540, USA and Institute for Advanced Study, Princeton, NJ 08540, USA
bsudakov@math.princeton.edu


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