

BOREL COMPLEXITY OF SETS OF NORMAL NUMBERS VIA GENERIC POINTS IN SUBSHIFTS WITH SPECIFICATION

DYLAN AIREY, STEVE JACKSON, DOMINIK KWIETNIAK, AND BILL MANCE

ABSTRACT. We study the Borel complexity of sets of normal numbers in several numeration systems. Taking a dynamical point of view, we offer a unified treatment for continued fraction expansions and base r expansions, and their various generalisations: generalised Lüroth series expansions and β -expansions. In fact, we consider subshifts over a countable alphabet generated by all possible expansions of numbers in $[0, 1)$. Then normal numbers correspond to generic points of shift-invariant measures. It turns out that for these subshifts the set of generic points for a shift-invariant probability measure is precisely at the third level of the Borel hierarchy (it is a $\mathbf{\Pi}_3^0$ -complete set, meaning that it is a countable intersection of F_σ -sets, but it is not possible to write it as a countable union of G_δ -sets). We also solve a problem of Sharkovsky–Sivak on the Borel complexity of the basin of statistical attraction. The crucial dynamical feature we need is a feeble form of specification. All expansions named above generate subshifts with this property. Hence the sets of normal numbers under consideration are $\mathbf{\Pi}_3^0$ -complete.

1. INTRODUCTION

Roughly speaking, a *numeration system* assigns to each real number an *expansion*. Here, an expansion is an infinite sequence of *digits* coming from some at most countable set. A real number is *normal* in a numeration system if all *asymptotic frequencies* of finite blocks of consecutive digits appearing in the expansion are *typical* for the numerations systems. To put some more content into this vague description recall that a real number ξ is normal in base 2 if in its binary expansion every block of digits of length k appears with asymptotic frequency $1/2^k$. It follows that for every integer $r \geq 2$ the set of normal numbers in base r is a first category set of full Lebesgue measure. Also, the normal numbers form a Borel set. As we explain below, the same holds true for all numeration systems we consider. For more on numeration systems, including different views on that theory see [?, ?, ?].

Knowing that the sets of normal numbers are Borel it is natural to gauge their complexity using the descriptive hierarchy of Borel sets. In that hierarchy, the simplest Borel sets are open ones and their complements (closed sets). On the next level, there are countable intersections and countable unions of sets at the first level. These are G_δ and F_σ sets, and the third level is formed by taking countable intersections and unions of sets at the second level. The procedure continues and provides a stratification of the family of Borel sets into levels corresponding to countable ordinals. It is known that for an uncountable Polish space these levels do not collapse: at each level there appear new sets which do not occur at any lower level of the hierarchy. Thus to every Borel set we can associate its complexity, that is, the lowest level of the hierarchy at which the set is visible. On the other hand, determining the position of “naturally arising” or “non-ad hoc” sets in the hierarchy

is a challenging problem. Only a small number of concrete examples are known to appear only above the third level.

A. Kechris asked in the 90's whether the set of real numbers that are normal in base two is an example of a Borel set properly located at the third level, which was later confirmed by H. Ki and T. Linton in [?]. More precisely, Ki and Linton showed that the set of numbers that are normal in an integer base $r \geq 2$ is a Π_3^0 -complete set, which means that this set is a countable intersection of F_σ sets and cannot be represented as a countable union of G_δ -sets. Since then many authors have studied the Borel complexity of various sets related to normal numbers, and have extended this result in various directions [?, ?, ?, ?].

Here we study analogous problems from the dynamical system perspective. It allows us to obtain a vast generalization of the Ki and Linton result. As our primary motivation are applications to numeration systems we restrict ourselves to symbolic dynamical systems (*subshifts* for short) and we will address more dynamical aspects of that theory in a forthcoming paper [?].

Before stating our main theorem, let us now briefly explain the connection between normal numbers and generic points for subshifts. If \mathcal{A} is a finite or countable¹ set, which we call the *alphabet*, then the *full shift space* over \mathcal{A} is the pair $(\mathcal{A}^\omega, \sigma)$ where \mathcal{A}^ω is endowed with the product topology induced by the discrete topology on \mathcal{A} , and σ stands for the shift map, which is given for $(x_n)_{n \in \omega} \in \mathcal{A}^\omega$ by $\sigma(x)_n = x_{n+1}$. By a *subshift* of \mathcal{A}^ω (or *over \mathcal{A}*) we mean a pair (X, σ) , where X is a nonempty closed shift-invariant subset of \mathcal{A}^ω , and σ is the shift map restricted to X . We also write \mathcal{A}^n for the set all of all *blocks of length n over \mathcal{A}* , that is, \mathcal{A}^n stands for the set of all finite sequences $w = w_1 \dots w_n$ with $w_j \in \mathcal{A}$ for $1 \leq j \leq n$ and $n \in \mathbb{N}$. As we will explain later, the set of sequences of digits which are expansions of real numbers defines a subshift for each of the numeration systems we consider. Furthermore, normal numbers in these numerations systems always correspond to generic points for some invariant measure of the associated subshift. Recall that a Borel probability measure μ on \mathcal{A}^ω is *shift-invariant* if $\mu(A) = \mu(\sigma^{-1}(A))$ for every Borel set $A \subseteq \mathcal{A}^\omega$. We say a shift-invariant measure μ is an invariant measure for a subshift X if X contains the support of μ , that is, $\mu(X) = 1$. An invariant measure μ is *ergodic* if for every Borel set $A \subseteq \mathcal{A}^\omega$ the condition $\sigma^{-1}(A) = A$ implies $\mu(A) \in \{0, 1\}$ (this is equivalent to saying that if $A \subseteq \sigma^{-1}(A)$ then $\mu(A) \in \{0, 1\}$). We say that a finite block $w \in \mathcal{A}^n$ *appears* in $x \in \mathcal{A}^\omega$ at the position $\ell \in \omega$ if $x_{\ell+i-1} = w_i$ for each $1 \leq i \leq n$. Let $e(w, x, N)$ be the number of times w appears in x at a position $\ell < N$. Let X be a subshift over \mathcal{A} and μ be an invariant measure. A point $x \in X$ is *generic*² for μ if for every finite block $w \in \mathcal{A}^n$ the set of positions at which w appears in x has the frequency equal to the measure of the set of all sequences starting with w , that is, if

$$\lim_{N \rightarrow \infty} \frac{e(w, x, N)}{N} = \mu([w]),$$

¹We call a set *countable* if it has the cardinality \aleph_0 , where \aleph_0 stands for the smallest infinite cardinal. We need this extra generality to cover continued fractions expansions and some generalised Lüroth series expansions.

²This definition differs from the usual definition of a generic point for Polish spaces (cf. [?, p. 1748]), but it is better adapted to the symbolic setting. The equivalence of these two definitions is easy to see (c.f. Corollary 18.3.11 of [?]).

where $[w] = \{z \in \mathcal{A}^\omega : z_0 = w_1, \dots, z_{n-1} = w_n\}$. By the shift-invariance of μ the measure of $[w]$ is equal to the μ -probability of the occurrence of w at any fixed position $\ell \in \omega$, that is,

$$\mu([w]) = \mu(\{z \in \mathcal{A}^\omega : z_\ell = w_1, \dots, z_{\ell+n-1} = w_n\}).$$

The ergodic theorem guarantees that for every shift-invariant ergodic measure μ the set of points generic for μ , denoted G_μ , has full measure (this is well-known for compact spaces, for the proof of this fact in the generality considered here, see [?, Lemma 2.2]). With this vocabulary the theorem of Ki and Linton becomes the statement that setting $X = \{0, 1, \dots, r-1\}^\omega$, the set of generic points for the Bernoulli measure μ (which is the product of the countable sequence of uniform probability measures on $\mathcal{A} = \{0, 1, \dots, r-1\}$) is a $\mathbf{\Pi}_3^0$ -complete set. It is then natural to ask for which subshifts (X, T) and measures μ one can prove a similar result about the Borel set $G_\mu \subseteq X$. In particular, we would like to know if the same result holds for other numeration systems than the classical base r -expansions. In terms of the theory of dynamical systems, this amounts to asking for which subshifts and invariant measures the Borel complexity of the set of generic points is a $\mathbf{\Pi}_3^0$ -complete set. Not surprisingly, we are not the first to pose this problem. When the present paper was being finished we learned that in the context of dynamical systems this question was first raised by A. Sharkovsky and his disciple A. Sivak (see [?], which quotes [?] and [?] as the primary sources, unfortunately these papers are not available in English). Sharkovsky and Sivak worked independently of the normal numbers community and used a slightly different language (for example, they called G_μ the *basin of attraction of μ*). Sharkovsky and Sivak noted that G_μ is always a Borel set lying at most at the third level of the hierarchy. It is also easy to see that G_μ may be empty if μ is not ergodic. Furthermore, there are easy examples with G_μ lying at the lowest level of the Borel hierarchy. To see that consider the unit circle $X = \mathbb{R}/\mathbb{Z}$, $\alpha \in \mathbb{R} \setminus \mathbb{Q}$, and let T act as $x \mapsto x + \alpha \bmod 1$. Then for every point $x \in \mathbb{R}/\mathbb{Z}$ its forward T -orbit is the sequence $\{n\alpha + x \bmod 1 : n \geq 0\}$, so each orbit is uniformly distributed mod 1, which means that every point in the circle is generic for the Lebesgue measure λ on \mathbb{R}/\mathbb{Z} , so $G_\lambda = \mathbb{R}/\mathbb{Z}$ is a clopen set. The same holds for *Sturmian subshifts*, which are symbolic dynamical models for irrational rotations of the circle (see [?, p. 321]). Sharkovsky and Sivak asked if their upper bound for the complexity of G_μ can be reached (see Problems 3 and 5 in [?]). As we noted above this asks for a Ki and Linton type result for dynamical systems.³ Because of the examples where G_μ is below the third level we see that some assumptions on the dynamical systems are required for such a result to hold. It turns out that it suffices to assume that the system has some form of the *specification property*. The original specification property was introduced by R. Bowen in his paper on Axiom A diffeomorphisms [?]. The specification property has played an important role in dynamics. We refer the reader to [?] for a discussion of the specification property and its many variants as well as their significance in dynamics. Our main result says that for a subshift (X, σ) possessing a feeble form of the specification property the set G_μ of generic points is $\mathbf{\Pi}_3^0$ -complete for every

³Note that the equivalence between normal numbers and generic points for the Bernoulli measure implies that the Ki and Linton result answers Problem 5 from [?] in the positive, but does not solve Problem 3 from that paper.

σ -invariant Borel probability measure μ . We also demonstrate that the theorem applies to many dynamical systems generating expansions of real numbers.

Thus the main theorem, which is to our best knowledge the first result of this type for dynamical systems, contains also several previously obtained results on complexity of sets of normal numbers, as well as many new ones. In particular, we extend the Ki-Linton result to continued fraction expansions, β -expansions, and generalized GLS⁴ expansions (of which the tent map is a special case).

In addition we note that there are subshifts, which are not so closely connected with numeration systems, but are interesting for the symbolic dynamics community, where our methods apply. These include hereditary subshifts (see Section 4 [?] for a more detailed overview).

In §2 we introduce basic definitions and notation, and mention the overall strategy. We introduce in this section the weak form of the specification property we require for our main result. In §3 we state and prove our main result. In §4 we give a number of applications of the main result including to continued fractions, β -expansions, generalized GLS-expansions. The enumeration system corresponding to the tent map is a special case of a generalized GLS expansion. This then answers a question of Sharkovsky–Sivak [?].

2. VOCABULARY/DEFINITIONS/NOTATION

Throughout this paper $\omega = \{0, 1, 2, \dots\}$ and $\mathbb{N} = \{1, 2, 3, \dots\}$. The cardinality of a finite set A is denoted by $|A|$. We write $\bar{d}(A)$ for the *upper asymptotic density* of a set $A \subseteq \omega$, that is,

$$\bar{d}(A) = \limsup_{n \rightarrow \infty} \frac{|(A \cap \{0, 1, \dots, n-1\})|}{n}.$$

2.1. Borel hierarchy. We now recall some basic notions from descriptive set theory which gauges the complexity of sets in Polish spaces. In any topological space X , the collection of Borel sets $\mathcal{B}(X)$ is the smallest σ -algebra containing all open sets. Elements of $\mathcal{B}(X)$ are stratified into levels, introducing the Borel hierarchy on $\mathcal{B}(X)$, by defining Σ_1^0 to be the family of open sets, and $\Pi_1^0 = \{X \setminus A : A \in \Sigma_1^0\}$ to be the family of closed sets. For a countable ordinal $\alpha < \omega_1$ we let Σ_α^0 be the collection of countable unions $A = \bigcup_n A_n$ where each $A_n \in \Pi_{\alpha_n}^0$ for some ordinal $\alpha_n < \alpha$. We also let $\Pi_\alpha^0 = \{X \setminus A : A \in \Sigma_\alpha^0\}$. Alternatively, $A \in \Pi_\alpha^0$ if $A = \bigcap_n A_n$ where $A_n \in \Sigma_{\alpha_n}^0$ and $\alpha_n < \alpha$ for each n . We also set $\Delta_\alpha^0 = \Sigma_\alpha^0 \cap \Pi_\alpha^0$ for each countable ordinal $\alpha < \omega_1$, in particular Δ_1^0 is the collection of clopen subsets of X . Note that Σ_2^0 is the collection of F_σ sets, and Π_2^0 is the collection of G_δ sets. For any topological space, $\mathcal{B}(X) = \bigcup_{\alpha < \omega_1} \Sigma_\alpha^0 = \bigcup_{\alpha < \omega_1} \Pi_\alpha^0$. It is easy to see that all of the collections Δ_α^0 , Σ_α^0 , Π_α^0 are *pointclasses*, that is, they are closed under inverse images of continuous functions. Another basic fact is that for any uncountable Polish space X , there is no collapse in the levels of the Borel hierarchy, that is, all the pointclasses Δ_α^0 , Σ_α^0 , Π_α^0 , for any ordinal $\alpha < \omega_1$, are distinct (for a proof, see [?]). Thus, these levels of the Borel hierarchy can be used to calibrate the descriptive complexity of a set. We say a set $A \subseteq X$ is Σ_α^0 (resp. Π_α^0) *hard* if $A \notin \Pi_\alpha^0$ (resp. $A \notin \Sigma_\alpha^0$). This says A is “no simpler” than a Σ_α^0 set. We say A is Σ_α^0 -*complete* if

⁴Note that GLS stands here for *generalized Lüroth series*, that is, the notion of a generalized GLS expansion is an extension of the GLS expansion, see [?] for more details.

$A \in \Sigma_\alpha^0 \setminus \Pi_\alpha^0$, that is, $A \in \Sigma_\alpha^0$ and A is Σ_α^0 hard. This says A is exactly at the complexity level Σ_α^0 . Likewise, A is Π_α^0 -complete if $A \in \Pi_\alpha^0 \setminus \Sigma_\alpha^0$.

Let us now discuss our proofs. In order to determine the exact position of a set A in the Borel hierarchy one must prove an *upper bound*, that is one must write a condition defining A which shows that it appears at some level in the hierarchy, and then show a *lower bound*, that is, show that A does not belong to any lower level in the hierarchy. To establish a lower bound we use a technique known as “Wadge reduction”. It is based on the observation that our hierarchy levels are all pointclasses. Thus, for example, a Borel set A is Σ_α^0 -hard if there are a Polish space Y , a Borel set $C \subseteq Y$ which is known to be Σ_α^0 -hard, and a continuous function $f: Y \rightarrow X$ such that $f^{-1}(A) = C$. The same holds for the Π_α^0 classes. Although the whole idea is plain and simple, the difficulty lies in the proper choice of the model space Y and subset C , so that it becomes possible to write down a definition of an appropriate continuous function.

2.2. Shift spaces. For a comprehensive introduction to symbolic dynamics we refer to the book [?] by Lind and Marcus. For a shift space $X \subseteq \mathcal{A}^\omega$ and integer $n \geq 1$, we write $\mathcal{L}_n(X) \subseteq \mathcal{A}^n$ for the set of n -blocks appearing in X , that is $w \in \mathcal{L}_n(X)$ if and only if there exists some $x \in X$ and $\ell \in \omega$ such that $x_{\ell+i-1} = w_i$ for all $1 \leq i \leq n$. The *length* of a block w over \mathcal{A} is the number of symbols in w and it is denoted by $|w|$. We agree that \mathcal{A}^0 consists of a single element, called the *empty word*, that is, \mathcal{A}^0 contains only the unique block over \mathcal{A} of length 0. By $\mathcal{A}^{<\omega}$ we denote the set of all finite blocks over \mathcal{A} (including the empty word). We let $\mathcal{L}(X) = \bigcup_{n \geq 1} \mathcal{L}_n(X)$ and call $\mathcal{L}(X)$ the language of X . Note that $\mathcal{L}(X)$ does not contain the empty word. For $n \geq 1$ and a block $w \in \mathcal{A}^n$, by $[w]$ we denote the cylinder consisting of those $x \in \mathcal{A}^\omega$ with $x_i = w_i$ for $1 \leq i \leq n$. If X is a subshift and $w \in \mathcal{L}_n(X)$, then we define $[w]_X = [w] \cap X$. When there is no ambiguity we drop the dependence on X in our notation and write just $[w]$ for $[w]_X$. Henceforth, we enumerate all nonempty blocks in $\mathcal{A}^{<\omega}$ in such a way that all blocks appear before their proper extensions. Any such enumeration induces an analogous enumeration on $\mathcal{L}(X)$ for every shift space $X \subseteq \mathcal{A}^\omega$, that is we can always write $\mathcal{L}(X) = \{w_1, w_2, \dots\}$ in such a way that if w_i is a proper initial segment of w_j , then $i < j$. In any such enumeration, we always have $|w_n| \leq n$ for every $n \geq 1$. Note that the whole theory of shift spaces remains the same if instead of \mathcal{A}^ω , we consider $\mathcal{A}^\mathbb{N}$.

2.3. Frequencies of subblocks. Recall that $e(w, x, N)$ denotes the number of times a block $w \in \mathcal{A}^{<\omega}$ appears in $x \in \mathcal{A}^\omega$ at a position $\ell < N$. Similarly, we write $e'(w, u)$ for the number of times w appears as a subblock of u . We agree that the empty word *never* appears as a subblock of a finite block. We say that a finite block u is (m, ε) -good for a shift-invariant measure μ if for every $1 \leq j \leq m$ the fraction of positions at which w_j appears as a subblock of u is ε -close to the μ -measure of the cylinder of w_j , that is, we have

$$(1) \quad \mu([w_j]) - \varepsilon < \frac{e'(w_j, u)}{|u|} < \mu([w_j]) + \varepsilon \quad \text{for } j = 1, \dots, m.$$

(Recall that we have fixed an enumeration of all blocks in $\mathcal{A}^{<\omega}$.) We say that a sequence $(u_n)_{n \in \mathbb{N}}$ of finite blocks $u_n \in \mathcal{A}^{<\omega}$ with $|u_n| \rightarrow \infty$ as $n \rightarrow \infty$ *generates* a

shift-invariant measure μ if for every $w \in \mathcal{A}^{<\omega}$ we have

$$\lim_{n \rightarrow \infty} \frac{e'(w, u_n)}{|u_n| - |w| + 1} = \mu([w]).$$

Equivalently, a sequence $(u_n)_{n \in \mathbb{N}}$ of elements of $\mathcal{A}^{<\omega}$ generates a shift-invariant measure μ if for every $m \in \mathbb{N}$ and $\varepsilon > 0$ there is an n_0 such that u_n is (m, ε) -good for μ for every $n \geq n_0$.

For $x \in \mathcal{A}^\omega$, $N \geq 1$, and $w \in \mathcal{A}^k$ we clearly have

$$(2) \quad e'(w, x_{[0, N]}) \leq e(w, x, N) \leq e'(w, x_{[0, N]}) + k - 1,$$

where $x_{[0, N]} = x_0 x_1 \dots x_{N-1}$. It follows that $x \in \mathcal{A}^\omega$ is a generic point for a shift-invariant measure μ if and only if the sequence $(x_{[0, N]})_{N \in \mathbb{N}}$ generates μ .

For further reference note that for every $u, v, w \in \mathcal{A}^{<\omega}$ the following holds

$$(3) \quad e'(w, v) \leq e'(w, u) + e'(w, v) \leq e'(w, uv) \leq e'(w, u) + e'(w, v) + |w| - 1.$$

Definition 1. A sequence $(u_n)_{n \in \mathbb{N}}$ of elements of $\mathcal{A}^{<\omega}$ is

- *dominating* if the sequence $(|u_1| + \dots + |u_n|)/|u_{n+1}|$ converges monotonically to 0 as $n \rightarrow \infty$,
- *asymptotically stable* for a shift-invariant measure μ if for every $\varepsilon > 0$ and $m \in \mathbb{N}$ there is $N \in \mathbb{N}$ such that for every $n > N$ there is some $\ell' < |u_n|$ so that $\ell'/|u_{n-1}| < \varepsilon$ and for every $\ell' \leq \ell \leq |u_n|$ the restriction of u_n to the first ℓ letters is (m, ε) -good for μ .

Lemma 2. *If a sequence $(u_n)_{n \in \mathbb{N}}$ of elements $\mathcal{A}^{<\omega}$ is dominating and asymptotically stable for a shift-invariant Borel probability measure μ , then $(u_n)_{n \in \mathbb{N}}$ generates μ and the point $x = u_1 u_2 u_3 \dots$ is generic for μ .*

Proof. It is clear that $|u_n| \rightarrow \infty$ as $n \rightarrow \infty$. The definition of asymptotic stability implies immediately that (u_n) generates μ . Let $U_n = u_1 u_2 \dots u_n$ for $n \geq 1$. Applying (3) to $U_n = U_{n-1} u_n$ we have for every $w \in \mathcal{A}^{<\omega}$ that

$$\frac{e'(w, u_n)}{|u_n|} \frac{|u_n|}{|U_n|} \leq \frac{e'(w, U_n)}{|U_n|} \leq \frac{|U_{n-1}|}{|u_n|} + \frac{e'(w, u_n)}{|u_n|} + \frac{|w| - 1}{|u_n|}.$$

Taking into account that the sequence (u_n) is dominating, so $|U_{n-1}|/|u_n|$ goes to 0 and $|U_n|/|u_n|$ converges to 1 as $n \rightarrow \infty$, we have for every $w \in \mathcal{A}^{<\omega}$ that

$$\lim_{n \rightarrow \infty} \frac{e'(w, U_n)}{|U_n|} = \lim_{n \rightarrow \infty} \frac{e'(w, u_n)}{|u_n|} = \mu([w]).$$

It remains to show that x is generic for μ . It is enough to show that for every $m \in \mathbb{N}$ and $\varepsilon > 0$ we can find $K > 0$ so that $x_{[0, k]}$ is (m, ε) -good for all $k \geq K$.

To this end fix $w \in \mathcal{A}^{<\omega}$ and consider the initial subblock $x_{[0, k]}$ of x . It follows that for all sufficiently large k we can write $x_{[0, k]} = U_n v$ for some $n \in \mathbb{N}$ and a proper subblock v of U_{n+1} . Pick $\varepsilon > 0$ and m large enough for w to be among w_1, \dots, w_m . Use m and $\varepsilon/2$ to find N as in the definition of asymptotic stability and assume that k is large enough so that the n for which $x_{[0, k]} = U_n v$ is strictly greater than N . For that n we can find ℓ' as in the definition of asymptotic stability. We have two cases to consider. First, if $|v| < \ell'$, then using (3) we get

$$e'(w, U_n) \leq e'(w, U_n v) \leq e'(w, U_n) + |v| + |w| - 1.$$

It follows that

$$(4) \quad \frac{e'(w, U_n)}{|U_n|} \frac{|U_n|}{|U_n v|} \leq \frac{e'(w, U_n v)}{|U_n v|} \leq \frac{e'(w, U_n)}{|U_n|} + \frac{\ell' + |w| - 1}{|U_n|}.$$

Since U_n is $(m, \varepsilon/2)$ -good for μ we can use (4) with (1) to get

$$(5) \quad (\mu([w]) - \varepsilon/2) \frac{|U_n|}{|U_n v|} \leq \frac{e'(w, U_n v)}{|U_n v|} \leq \mu([w]) + \varepsilon/2 + \frac{\ell' + |w| - 1}{|U_n|}.$$

Now the left hand side of (5) satisfies

$$(\mu([w]) - \varepsilon/2) \frac{|U_n|}{|U_n v|} \geq \mu([w]) \left(1 - \left(1 - \frac{|U_n|}{|U_n v|}\right)\right) - \varepsilon/2 \geq \mu([w]) - \varepsilon/2 - \frac{|v|}{|U_n v|}.$$

Plugging that into (5) we obtain

$$(6) \quad \mu([w]) - \varepsilon/2 - \frac{|v|}{|U_n v|} \leq \frac{e'(w, U_n v)}{|U_n v|} \leq \mu([w]) + \varepsilon/2 + \frac{\ell' + |w| - 1}{|U_n|}.$$

In the second case $|v| \geq \ell'$, which implies that v is $(m, \varepsilon/2)$ -good for μ . By (3) we obtain

$$(7) \quad e'(w, U_n) + e'(w, v) \leq e'(w, U_n v) \leq e'(w, U_n) + e'(w, v) + |w| - 1.$$

Being $(m, \varepsilon/2)$ -good for μ (see (1)) means that

$$(8) \quad \mu([w])|U_n| - \varepsilon|U_n|/2 < e'(w, U_n) < \mu([w])|U_n| + \varepsilon|U_n|/2$$

and

$$(9) \quad \mu([w])|v| - \varepsilon|v|/2 < e'(w, v) < \mu([w])|v| + \varepsilon|v|/2.$$

Applying (8) and (9) to (7) we obtain that

$$(10) \quad \mu([w]) - \varepsilon/2 < \frac{e'(w, U_n v)}{|U_n| + |v|} < \mu([w]) + \varepsilon/2 + \frac{|w| - 1}{|U_n| + |v|}.$$

Now, (6) and (10) imply that for all sufficiently large n the block $U_n v$ is (m, ε) -good for μ . \square

Let d_H stand for the normalised Hamming distance, that is, given two blocks $u = u_1 \dots u_n$ and $w = w_1 \dots w_n$ of equal length we set $d_H(u, w) = |\{1 \leq j \leq n : u_j \neq w_j\}|/n$.

Lemma 3. *Suppose $x, y \in \mathcal{A}^\omega$ and $x \in G_\mu$ for a shift-invariant Borel probability measure μ on \mathcal{A}^ω .*

- (a) *If $\bar{d}(x, y) = \bar{d}(\{j \in \omega : x_j \neq y_j\}) = 0$, then $y \in G_\mu$.*
- (b) *If $y = x_{i_0} x_{i_1} x_{i_2} \dots$ where $(i_j)_{j \in \omega}$ is a strictly increasing sequence of elements of ω such that $\bar{d}(\omega \setminus \{i_j : j \in \omega\}) = 0$, then $y \in G_\mu$ iff $x \in G_\mu$.*
- (c) *Let $x = u_1 u_2 u_3 \dots$, $y = v_1 v_2 v_3 \dots$, where $(u_n)_{n \in \mathbb{N}}$ and $(v_n)_{n \in \mathbb{N}}$ are sequences of blocks in $\mathcal{A}^{<\omega}$ satisfying $|u_n| = |v_n|$ for every $n \geq 1$. If $d_H(u_n, v_n) \rightarrow 0$ as $n \rightarrow \infty$, and*

$$(11) \quad \frac{|u_{n+1}| \cdot d_H(u_{n+1}, v_{n+1})}{|u_1| + \dots + |u_n|} \rightarrow 0 \quad \text{as } n \rightarrow \infty,$$

then $\bar{d}(x, y) = \bar{d}(\{j \in \omega : x_j \neq y_j\}) = 0$.

- (d) *For every $m \in \mathbb{N}$ and $\varepsilon > 0$ there exists $\delta > 0$ such that if $w \in \mathcal{A}^n$ is $(m, \varepsilon/2)$ -good and $w' \in \mathcal{A}^n$ satisfies $d_H(w, w') < \delta$, then w' is (m, ε) -good.*

Proof. The first two statements can be found in [?, ?]. The proof of the fourth is straightforward. Let us prove the third one. Fix $\varepsilon > 0$. Let $\ell_n = |u_1| + \dots + |u_n|$ for $n \in \mathbb{N}$. Since $d_H(u_n, v_n) \rightarrow 0$ as $n \rightarrow \infty$, we find K such that

$$(12) \quad d_H(u_n, v_n) < \varepsilon \quad \text{for all } n \geq K.$$

Using (11) we obtain $N > K$ such that for every $n \geq N$ we have

$$(13) \quad |u_{n+1}| \cdot d_H(u_{n+1}, v_{n+1}) < \varepsilon \cdot \ell_n.$$

We can pick N so large that $\ell_K/\ell_N < \varepsilon$ holds. Then for $n > N$ and every j such that $\ell_n \leq j < \ell_{n+1}$ it holds that

$$\begin{aligned} \frac{|\{0 \leq i < j : x_i \neq y_i\}|}{j} &\leq \frac{1}{\ell_n} \cdot \left(\sum_{k=1}^{n+1} |u_k| \cdot d_H(u_k, v_k) \right) \\ &\leq \frac{\sum_{k=1}^K |u_k|}{\ell_n} + \frac{\sum_{k=K+1}^n |u_k| \cdot d_H(u_k, v_k)}{\ell_n} + \frac{|u_{n+1}| \cdot d_H(u_{n+1}, v_{n+1})}{\ell_n}. \end{aligned}$$

Now, note that the first term in the sum above equals ℓ_K/ℓ_n , and so is bounded by $\ell_K/\ell_N < \varepsilon$. By (13), the third term in the sum above is also bounded by ε . Finally, the same holds for the middle term, because $|u_{K+1}| + \dots + |u_n| \leq \ell_n$ and (12). We have proved that for every $\varepsilon > 0$ and all sufficiently large j it holds that

$$\frac{|\{0 \leq i < j : x_i \neq y_i\}|}{j} \leq 3\varepsilon,$$

which is what we wanted to show. \square

2.4. Specification for subshifts. For the general definition of the specification property we refer the reader to [?]. We omit it here, as for shift spaces it has a simple combinatorial reformulation. The equivalence of these two definitions is an easy exercise.

Definition 4. A shift space X over an at most countable alphabet \mathcal{A} has the *specification property* if there is a nonnegative integer N such that if $w_i \in \mathcal{L}(X)$ for $i = 1, \dots, n$ then there are $v_i \in \mathcal{A}^N$ for $i = 1, \dots, n-1$ such that $u = w_1 v_1 w_2 v_2 \dots v_{n-1} w_n \in \mathcal{L}(X)$. Furthermore, we say that X has the *periodic specification property* if, in addition to $v_i \in \mathcal{A}^N$ for $i = 1, \dots, n-1$ as above we can take v_n so that the periodic point $x = (w_1 v_1 w_2 v_2 \dots w_n v_n)^\infty$ belongs to X .

Note that if X is a compact subshift, then the specification property and its periodic version are well known to be equivalent. Also, when X is not compact then the specification property is no longer necessarily an invariant for topological conjugacy.

The classical specification property is much too strong for our purposes as it does not apply to most β -shifts. It is then natural to replace it by a weaker assumption. Looking for such a notion we found out that no existing generalisation of the specification property is fully satisfactory. Therefore we introduce yet another property, which we coin the *right feeble specification property*. It is similar to the almost specification property, which was originally defined by Pfister and Sullivan [?], and later modified and renamed by Thompson [?]. The reader may consult [?] for discussion of this property. A variant of the latter property, the *right almost specification property*, was considered by Climenhaga and Pavlov (for more details we refer the reader to Definition 2.14 in [?]). We need a similar condition here to

guarantee that the function we will define in the course of our proof of Theorem 6 is continuous.

Definition 5. We say that a subshift X has the *right feeble specification* property if there exists a set $\mathcal{G} \subseteq \mathcal{L}(X)$ satisfying:

- (1) a concatenation of words in \mathcal{G} stays in \mathcal{G} , that is, if $u, v \in \mathcal{G}$, then $uv \in \mathcal{G}$;
- (2) for any $\epsilon > 0$ there is an $N = N(\epsilon)$ such that for every $u \in \mathcal{G}$ and $v \in \mathcal{L}(X)$ with $|v| \geq N$, there are $s, v' \in \mathcal{A}^{<\omega}$ satisfying $|v'| = |v|$, $0 \leq |s| \leq \epsilon|v|$, $d_H(v, v') < \epsilon$, and $usv' \in \mathcal{G}$.

It is immediate that the right almost specification property implies the right feeble specification, in particular the specification property implies the feeble specification property (cf. [?, Lemma 2.15]). It is also easy to see that the weak specification property (see [?, ?]) implies the right feeble specification. We do not know if the weak specification property (or the right feeble specification property) implies the right almost specification property. We suspect that the answer to both questions is “no” and an appropriate example can be constructed within the family of subshifts with the weak specification property presented in [?].

2.5. Irregular set. Given $w \in \mathcal{L}(X)$ we define $I_w(X)$ to be the set of all $x \in X$ such that the set of positions at which w appears in x does not have a frequency, that is

$$\liminf_{N \rightarrow \infty} \frac{e(w, x, N)}{N} < \limsup_{N \rightarrow \infty} \frac{e(w, x, N)}{N}.$$

Let $I(X)$ be the *irregular set* for X , that is, the union of sets $I_w(X)$ over all $w \in \mathcal{L}(X)$. The *quasi-regular set* for X is the complement of $I(X)$, that is, $Q(X) = X \setminus I(X)$. Both sets are obviously Borel and belong to the third level of the Borel hierarchy.

3. MAIN RESULTS

3.1. Subshifts with a feeble specification property. Theorem 6 below applies to subshifts on a countable alphabet satisfying a hypothesis weaker than the (non-periodic) specification property.

Note that we are considering subshifts which are not necessarily compact. It forces us to *assume* that there are at least two shift-invariant Borel probability measures on X . This condition is automatically fulfilled if X is compact with $|X| \geq 2$ (granting the specification hypothesis).

Theorem 6. Assume that \mathcal{A} is at most countable and X is a subshift over \mathcal{A} with the right feeble specification property and at least two shift-invariant Borel probability measures. If μ is a shift-invariant Borel probability measure on X , then every Borel set B satisfying $G_\mu \subseteq B \subseteq Q(X)$, where G_μ is the set of generic points for μ and $Q(X)$ is the quasi-regular set, is Π_3^0 -hard. In particular, B is Π_3^0 -complete provided that B is a Π_3^0 -set. Hence, G_μ and $Q(X)$ are Π_3^0 -complete, and the irregular set $I(X)$ is Σ_3^0 -complete.

Proof. First, we note that under our assumptions the set of generic points is nonempty for every shift-invariant Borel probability measure on X . The existence of a generic point for an *arbitrary* shift-invariant measure follows from the right

feeble specification property and Corollary 22 in [?] (formally, the quoted result requires a stronger assumption, but the proof remains the same when we just assume right feeble specification property).

Fix a shift-invariant Borel probability measure μ on X and a Borel set B such that $G_\mu \subseteq B \subseteq Q(X)$. Let $\varepsilon \mapsto N(\varepsilon)$ be the function as implicitly defined for X by Definition 5. In order to apply Wadge reduction, it suffices to find a Polish metric space \mathcal{X} , a continuous function $\pi: \mathcal{X} \rightarrow X$ and a $\mathbf{\Pi}_3^0$ -complete set $\mathcal{C}_3 \subseteq \mathcal{X}$ such that

$$\pi^{-1}(G_\mu) = \pi^{-1}(B) = \pi^{-1}(Q(X)) = \mathcal{C}_3$$

and $\pi^{-1}(I(X)) = \mathcal{X} \setminus \mathcal{C}_3$.

We take $\mathcal{X} = \mathbb{N}^{\mathbb{N}}$ with the topology of pointwise convergence, and choose $\mathcal{C}_3 \subseteq \mathbb{N}^{\mathbb{N}}$ to be the set of all functions $\alpha: \mathbb{N} \rightarrow \mathbb{N}$ attaining any $n \in \mathbb{N}$ only finitely many times, that is,

$$\mathcal{C}_3 = \{\alpha \in \mathbb{N}^{\mathbb{N}}: \liminf_{n \rightarrow \infty} \alpha(n) = \infty\}.$$

It is well-known that \mathcal{C}_3 is a $\mathbf{\Pi}_3^0$ -complete set.

In order to define π we fix a shift-invariant Borel probability measure $\nu \neq \mu$ on X . Then we fix a μ -generic point $x \in X$ and a ν generic point $z \in X$.

We will also need auxiliary ω -valued sequences $(a_n)_{n \geq 0}$, $(b_n)_{n \geq 0}$, and $(c_n)_{n \geq 0}$ to be defined in a moment. It is also convenient to introduce one more auxiliary sequence $(B_n)_{n \geq 0}$, so that $B_0 = 0$ and $B_k = 2(b_1 + \dots + b_k)$ for $k \geq 1$. Given $\alpha \in \mathbb{N}^{\mathbb{N}}$ and using these sequences we define blocks $u_1, u_2, \dots \in \mathcal{L}(X)$ inductively, defining a group of cardinality $2b_n$ at one step, first for $1 \leq j \leq 2b_1$, by

$$(14) \quad u_j = \begin{cases} x_{[0, a_1)}, & \text{if } 0 < j \leq b_1, \text{ and} \\ z_{[0, c_1)}, & \text{if } b_1 < j \leq 2b_1, \end{cases}$$

and then, assuming that u_1, \dots, u_i have been defined where $i = 2(b_1 + \dots + b_n) = B_n$ for some $n \geq 1$, we set

$$(15) \quad u_j = \begin{cases} x_{[0, a_{n+1})}, & \text{if } B_n < j \leq B_n + b_{n+1}, \text{ and} \\ z_{[0, c_{n+1})}, & \text{if } B_n + b_{n+1} < j \leq B_n + 2b_{n+1}. \end{cases}$$

We now want to produce finite blocks $v_0, v_1, v_2, v_3, \dots$ in $\mathcal{L}(X)$ so that all the concatenations $v_0 v_1 v_2 \dots v_n$ for $n \geq 1$ are in $\mathcal{L}(X)$ and for each $j \geq 1$ the block v_j is close (in an appropriate sense) to u_j .

To do so we apply the right feeble specification property inductively. We start with an arbitrarily chosen $v_0 \in \mathcal{G}$. Assume that we have defined v_1, \dots, v_{j-1} for some $j \geq 1$. Then we use the right feeble specification property to obtain v_j so that $v_1 v_2 \dots v_j \in \mathcal{G}$ and we have $v_j = s_j u'_j$ where u'_j has the same length as u_j , the length of s_j is a tiny fraction of $|u_j|$, and the Hamming distance $d_H(u'_j, u_j)$ is small.

We will then set $\pi(\alpha) = \sigma^{|v_0|}(v_0 v_1 v_2 \dots) = v_1 v_2 v_3 \dots$. Note that $\pi(\alpha) \in X$ because X is closed and shift invariant. With the right choice of the a_n 's, b_n 's, and c_n 's we will prove that the map $\alpha \mapsto \pi(\alpha)$ is the required reduction.

Now we will define our auxiliary sequences. For $\alpha \in \mathbb{N}^{\mathbb{N}}$, let $\alpha'(n) = \min\{n, \alpha(n)\}$. Let $(a_n)_{n \geq 0}$, $(c_n)_{n \geq 0}$ be sequences of positive integers with $a_0 = c_0 = 1$ growing so fast that for every $n \in \mathbb{N}$ the following conditions hold:

$$(16a) \quad a_n = \alpha'(n) c_n,$$

$$(16b) \quad c_n / n > 2^{2^n},$$

$$(16c) \quad c_n > N(1/2^{2n}),$$

$$(16d) \quad \text{for each } m \geq c_n \text{ the block } x_{[0,m)} \text{ is } (m, 1/2^{n+1})\text{-good for } \mu, \text{ and}$$

$$(16e) \quad \text{for each } m \geq c_n \text{ the block } z_{[0,m)} \text{ is } (m, 1/2^{n+1})\text{-good for } \nu.$$

Now define $b_0 = 0$ and $(b_n)_{n \geq 1}$ to be a sequence of positive integers satisfying for every $n \geq 1$ the following conditions:

$$(17a) \quad b_n > 2^{2n},$$

$$(17b) \quad a_n b_n > 2^{2n} a_{n+1}, \text{ and}$$

$$(17c) \quad a_n b_n > 2^{2n}((a_1 + c_1)b_1 + \dots + (a_{n-1} + c_{n-1})b_{n-1}).$$

Equations (14) and (15) now define the blocks u_n for $n \geq 1$. For $n \geq 1$ let

$$\begin{aligned} \bar{u}'_n &= (x_{[0,a_n)})^{b_n} = \underbrace{x_{[0,a_n)} \cdots x_{[0,a_n)}}_{b_n \text{ times}}, \text{ and} \\ \bar{u}''_n &= (z_{[0,c_n)})^{b_n} = \underbrace{z_{[0,c_n)} \cdots z_{[0,c_n)}}_{b_n \text{ times}}. \end{aligned}$$

Note that \bar{u}'_n is the concatenation of the u_i 's where i runs from $B_{n-1}+1$ to $B_{n-1}+b_n$, and \bar{u}''_n is the concatenation of the u_i 's where i runs from $B_{n-1}+b_n+1$ to $B_n = B_{n-1}+2b_n$ for each $n \geq 1$. It follows that the points $u_1 u_2 u_3 \dots \in \mathcal{A}^\omega$ and $\bar{u}'_1 \bar{u}''_1 \bar{u}'_2 \bar{u}''_2 \dots$ are equal.

We claim that:

- (A) If $\alpha \in \mathcal{C}_3$, then $(\bar{u}'_n \bar{u}''_n)$ is dominating and an asymptotically stable sequence for μ , which implies by Lemma 2 that the point $x = u_1 u_2 u_3 \dots$ is generic for μ .
- (B) If $\alpha \notin \mathcal{C}_3$, then the sequence (U'_n) given by

$$U'_n = \bar{u}'_1 \bar{u}''_1 \bar{u}'_2 \bar{u}''_2 \dots \bar{u}'_{n-1} \bar{u}''_{n-1} \bar{u}'_n$$

generates μ and the sequence (U''_n) given by

$$U''_n = \bar{u}'_1 \bar{u}''_1 \bar{u}'_2 \bar{u}''_2 \dots \bar{u}'_n \bar{u}''_n = U'_n \bar{u}''_n$$

generates along some subsequence a measure ν' , which is a (nontrivial) convex combination of μ and ν . Since $\nu' \neq \mu$, we see that $x = u_1 u_2 u_3 \dots$ is an irregular point.

Proof of Claim (A). Assume that $\alpha \in \mathcal{C}_3$. We first prove the following claim:

- (A') For each $m \in \mathbb{N}$ and $\varepsilon > 0$ the first ℓ symbols of the block \bar{u}'_n are (m, ε) -good for every sufficiently large $n \geq m$ and $\ell \geq \ell' = 2^n a_n$.

To see that take any $n \geq m$ and recall that $\bar{u}'_n = (x_{[0,a_n)})^{b_n}$ and $b_n > 2^n$ by (17a). Hence we can consider $2^n a_n = \ell' \leq \ell < |\bar{u}'_n| = a_n b_n$ and write \bar{u}'_n restricted to the first ℓ symbols as a concatenation $\tilde{u}'_n \tilde{u}''_n$ where $\tilde{u}'_n = (x_{[0,a_n)})^r$, $r = \lfloor \ell/a_n \rfloor \geq 2^n$ and $|\tilde{u}''_n| < a_n$. We have for every $1 \leq j \leq m$ that

$$ra_n e'(w_j, x_{[0,a_n)}) \leq e'(w_j, \tilde{u}'_n \tilde{u}''_n) \leq ra_n e'(w_j, x_{[0,a_n)}) + |\tilde{u}''_n| + |w_j| - 1.$$

Note that $|\tilde{u}''_n| + |w_j| \leq (1/2^n)ra_n + m$ and $m/a_n < m/c_n < 1/2^{2n}$ by (16a) and (16b). Now reasoning as in the proof of Lemma 2 and using the fact (16d) that $x_{[0,a_n)}$ is $(m, \varepsilon/2)$ -good for μ for all sufficiently large n we see that Claim (A') holds. To finish the proof of Claim (A) note that

$$|\bar{U}_n| = |\bar{u}'_n \bar{u}''_n| = (a_n + c_n)b_n = (\alpha'(n) + 1)c_n b_n.$$

Therefore (17c) implies that $(\bar{U}_n)_{n \in \mathbb{N}}$ is a dominating sequence, and Claim (A') together with (17b) and the fact that $\alpha'(n) \rightarrow \infty$ as $n \rightarrow \infty$ imply that $(\bar{U}_n)_{n \in \mathbb{N}}$ is an asymptotically stable sequence, so we can apply Lemma 2. (Claim (A) ■)

Proof of Claim (B). Observe that Claim (A') and (17c) imply that the sequence $(U'_n)_{n \in \mathbb{N}}$, where $U'_n = \bar{u}'_1 \bar{u}''_1 \bar{u}'_2 \bar{u}''_2 \dots \bar{u}'_{n-1} \bar{u}''_{n-1} \bar{u}'_n$, generates μ . We also have that there exists $M \in \mathbb{N}$ so that $\alpha'(n) = M$ for infinitely many n 's. Passing to a subsequence we can assume that this happens for all n . The same reasoning as in the proof of Claim (A') with (16e) replacing (16d) yields that the sequence $(U''_n)_{n \in \mathbb{N}}$, where $U''_n = \bar{u}'_1 \bar{u}''_1 \bar{u}'_2 \bar{u}''_2 \dots \bar{u}'_n \bar{u}''_n = U'_n \bar{u}''_n$, generates the measure $(1/(M+1)\nu + M/(M+1)\mu)$, which implies that $x = u_1 u_2 u_3 \dots$ is an irregular point. (Claim (B) ■)

Unfortunately, we cannot take $\pi(\alpha) = x = u_1 u_2 u_3 \dots$ because x need not belong to X . But given x we can use the right feeble specification property to find the sequence of blocks v_0, v_1, v_2, \dots as outlined above so that $\pi(\alpha) = \sigma^{|v_0|}(v_0 v_1 v_2 \dots) = v_1 v_2 v_3 \dots \in X$ and our construction will allow us to use Lemma 3 to prove that $\pi(\alpha)$ behaves like x .

We start with an arbitrarily chosen $v_0 \in \mathcal{G}$. Next we find v_1 such that $v_0 v_1 \in \mathcal{G}$ and $v_1 = s_1 u'_1$ where $|u'_1| = |u_1|$, $|s_1| < |u_1|/4$, and the Hamming distance $d_H(u_1, u'_1)$ is small (say, $d_H(u_1, u'_1) < 1/4$). Note that using Lemma 3(d) and inequality (3) (and increasing a_1 if necessary) we can assume that $v_1 \in \mathcal{G}$ is almost as good for μ as u_1 . Assume $v_0, v_1, v_2, \dots, v_i$ have been defined for some $i \geq 1$ so that $v_0 v_1 v_2 \dots v_i \in \mathcal{G}$ and $B_{n-1} \leq i < B_{n-1} + 2b_n$. Then the right feeble specification property gives us blocks s_{i+1} and u'_{i+1} such that

$$|u'_{i+1}| = |u_{i+1}| = \begin{cases} |x_{[0, a_n]}| = a_n, & \text{if } i < B_{n-1} + b_n, \text{ and} \\ |z_{[0, c_n]}| = c_n, & \text{if } B_{n-1} + b_n \leq i, \end{cases}$$

and $|s_{i+1}| \leq (1/2^{2n}) \cdot |u_i|$ and $d_H(u_{i+1}, u'_{i+1}) < 1/2^{2n}$ (we use here (16a) and (16c)). We set $v_{i+1} = s_{i+1} u'_{i+1}$ and let $\pi(\alpha) = \sigma^{|v_0|}(v_0 v_1 v_2 \dots) = v_1 v_2 v_3 \dots$. Note that the right feeble specification property guarantees that $v_0 v_1 v_2 \dots v_i v_{i+1} \in \mathcal{G}$. Define

$$I' = \{j \in \omega : j < |s_1| \text{ or } \exists n \geq 1 \text{ and } 0 < i \leq |s_{n+1}| \text{ with } j = |v_1 \dots v_n| + i\}.$$

Since $|s_{n+1}|/|v_{n+1}|$ goes to 0 as $n \rightarrow \infty$ we get that $\bar{d}(I') = 0$. Therefore Lemma 3(b) implies that $\pi(\alpha)$ is generic for μ if and only if $y = u'_1 u'_2 u'_3 \dots \in G_\mu$. Putting $x = u_1 u_2 u_3 \dots$ and $y = u'_1 u'_2 u'_3 \dots$ into Lemma 3(c) we see that $y \in G_\mu$ if and only if $x \in G_\mu$. Note that (11) holds because for every $n \geq 2$ and $B_{n-1} \leq i < B_{n-1} + 2b_n$ we have

$$\frac{|u_{i+1}| \cdot d_H(u_{i+1}, u'_{i+1})}{|u_1 \dots u_i|} \leq \frac{a_n \cdot (1/2^{2n})}{|U''_{n-1}|} \leq \frac{a_{n-1} b_{n-1} \cdot (1/2^{4n-2})}{(a_1 + c_1)b_1 + \dots + (a_{n-1} + c_{n-1})b_{n-1}},$$

which clearly implies (11) (for the last inequality we have used (17b) to bound the numerator). Similarly, we obtain that $\pi(\alpha) \in I(X)$ if and only if $y \in I(\mathcal{A}^\omega)$ if and only if $x \in I(\mathcal{A}^\omega)$. We conclude $\pi^{-1}(G_\mu) = \pi^{-1}(Q(X)) = \mathcal{C}_3$, and $\pi^{-1}(I(X)) = \mathbb{N}^\mathbb{N} \setminus \mathcal{C}_3$. These observations together with Claims (A) and (B) prove that the map $\alpha \mapsto \pi(\alpha)$ is a reduction map showing that B is Π_3^0 -hard, and so in particular G_μ and $Q(X)$ are Π_3^0 -complete and $I(X)$ is Σ_3^0 -complete, provided the map π is continuous. But the continuity is obvious as each initial segment of $\pi(\alpha)$ depends only on $\alpha(1), \dots, \alpha(n)$ for some $n \in \mathbb{N}$. □

Remark 7. Theorem 6 holds for any shift-invariant G_δ subset of \mathcal{A}^ω with the periodic specification property. The proof requires only a minor modification which we leave for the reader.

3.2. Hereditary Shifts. In §4 we present a number of applications of Theorem 6 to normal numbers defined by using various expansions including β -expansions, regular continued fraction expansions, and generalized Lüroth series expansions. In the remainder of this section we consider a result which does not follow immediately as a corollary to Theorem 6, but whose proof uses the same techniques as the one for that theorem. Namely, we show that the conclusion of Theorem 6 holds for the class of *hereditary shifts*. Furthermore, we can use Theorem 10 instead of Theorem 6 in the applications presented in §4, because every subshift considered there is hereditary and has a generic point for each of its invariant measures⁵. Actually, Theorem 10 is valid for an even broader class of subshifts having a *safe symbol* (see [?]).

Hereditary subshifts were introduced by Kerr and Li in [?, p. 882] (see also [?]). The family of hereditary subshifts includes extensively studied classes of subshifts: spacing shifts, β -shifts, bounded density shifts, \mathcal{B} -admissible shifts; also, many examples of \mathcal{B} -free shifts. Note also that all full shifts over $\{0, 1, \dots, n\}$ or ω are hereditary, as well as many sofic shifts and shifts of finite type (golden mean shift for example) (see Section 4 in [?] for more details and references).

Definition 8. A subshift $X \subseteq \mathcal{A}^\omega$ where $\mathcal{A} = \{0, 1, \dots, n\}$ or $\mathcal{A} = \omega$ is *hereditary* if $y \leq x$ coordinate-wise and $x \in X$ imply $y \in X$.

Definition 9. We say that $\gamma \in \mathcal{A}$ is a *safe symbol* for a subshift X over \mathcal{A} if for every $x \in X$ and $k \geq 0$ we have that the point y , where

$$y_n = \begin{cases} x_n, & \text{if } n \neq k, \text{ and} \\ \gamma, & \text{if } n = k, \end{cases}$$

also belongs to X .

Note that by definition 0 is a safe symbol for every hereditary subshift, and a subshift over $\{0, 1\}$ is hereditary if and only if 0 is its safe symbol. It is easy to see examples of subshifts over $\{0, 1, 2\}$ which are not hereditary but have 0 as a safe symbol. Shifts with a safe symbol seem to be more important in higher dimensional symbolic dynamics, see [?] and references therein. Note that we again need to *assume* that there are at least two shift-invariant measures on X , as even compact hereditary shifts may have only one invariant measure.

Theorem 10. Assume that \mathcal{A} is at most countable and X is a subshift over \mathcal{A} with a safe symbol (in particular, if $\mathcal{A} = \{0, 1, \dots, n\}$ or $\mathcal{A} = \omega$ and X is a hereditary shift) with more than one shift-invariant Borel probability measure. If μ is a shift-invariant Borel probability measure on X such that $G_\mu \neq \emptyset$ and B is a Borel set satisfying $G_\mu \subseteq B \subseteq Q(X)$, then B is Π_3^0 -hard. In particular, B is Π_3^0 -complete provided that B is a Π_3^0 -set. Hence, the set G_μ is either empty, or is Π_3^0 -complete. In particular, the set G_μ is Π_3^0 -complete for every ergodic measure. Furthermore, $Q(X)$ is a Π_3^0 -complete set and $I(X)$ is a Σ_3^0 -complete set.

⁵But the proof of the latter fact is anyway based on the specification property.

Proof. As in the proof of Theorem 6, we are going to define a continuous reduction $\pi: \mathbb{N}^{\mathbb{N}} \rightarrow X$ with $\pi^{-1}(G_\mu) = \pi^{-1}(Q(X)) = \mathcal{C}_3 = \{\beta \in \mathbb{N}^{\mathbb{N}}: \liminf_{n \rightarrow \infty} \beta(n) = \infty\}$ and $\pi^{-1}(I(X)) = \mathbb{N}^{\mathbb{N}} \setminus \mathcal{C}_3$. Without loss of generality we assume that 0 is a safe symbol for X . By δ_0 we denote the Dirac measure concentrated on $0^\infty = 000 \dots \in X$. Let μ be any shift-invariant measure on X . Suppose first that $\mu \neq \delta_0$, that is, μ is not supported on $\{0^\infty\}$. Then $\mu([\gamma]) > 0$ for some $\gamma \in \{1, \dots, n-1\}$ by the invariance of μ . Assume that G_μ is nonempty and take $x \in G_\mu$. Fix a strictly increasing sequence of nonnegative integers (b_n) such that $b_0 = 0$,

$$(18) \quad \lim_{n \rightarrow \infty} \frac{b_n}{b_{n+1}} = 0, \text{ and}$$

$$(19) \quad \lim_{n \rightarrow \infty} \frac{|\{k \in [b_n, b_{n+1}): x_k = \gamma\}|}{b_{n+1} - b_n} = \mu([\gamma]).$$

Fix $\beta \in \mathbb{N}^{\mathbb{N}}$. Let $n \in \mathbb{N}$ and let I_n be the set of positions in $[b_{2n-1}, b_{2n})$ where γ appears in x , that is,

$$I_n = \{k \in \mathbb{N} : b_{2n-1} \leq k < b_{2n} \text{ and } x_k = \gamma\}.$$

Let $q_n = |I_n|$. Write $I_n = \{i_1, \dots, i_{q_n}\}$ where $b_{2n-1} \leq i_1 < i_2 < \dots < i_{q_n} < b_{2n}$. Let $p_n = q_n - \lfloor q_n/\beta(n) \rfloor + 1$ and $J_n = \{i_{p_n}, i_{p_n+1}, \dots, i_{q_n}\}$. Note that $q_n/\beta(n) \leq |J_n| = \lceil q(n)/\beta(n) \rceil \leq q_n/\beta(n) + 1$. Define $\pi: \mathbb{N}^{\mathbb{N}} \rightarrow X$ by $\pi(\beta) = y$ where

$$y_k = \begin{cases} 0, & \text{if } k \in \bigcup J_n, \text{ and} \\ x_k, & \text{otherwise.} \end{cases}$$

Note that we have defined y so that it agrees with x except on the positions in the set

$$\bigcup_{n \in \mathbb{N}} J_n \subseteq \bigcup_{n \in \mathbb{N}} [b_{2n-1}, b_{2n}).$$

In particular, for each $n \geq 0$ we have

$$(20) \quad x_{[b_{2n}, b_{2n+1})} = y_{[b_{2n}, b_{2n+1})}.$$

Note also that for each $n \in \mathbb{N}$ to get $y_{[0, b_{2n})}$ we modify $x_{[0, b_{2n})}$ along at most

$$(21) \quad b_{2n-1} + \left(\frac{b_{2n} - b_{2n-1}}{\beta(n)} \right)$$

positions. We have $y \in X$ for every $\beta \in \mathbb{N}^{\mathbb{N}}$ since $y = \pi(\beta)$ is obtained from $x \in X$ by setting x_k to 0 for $k \in \bigcup J_n$ and 0 is the safe symbol for X . The map π is continuous since for each $n \in \mathbb{N}$ it is easy to see that $y_{[0, b_{2n})}$ depends only on x and $\beta(1), \dots, \beta(n)$.

If $\beta \in \mathcal{C}_3$ then $\lim_{n \rightarrow \infty} \beta(n) = \infty$ so the set $\bigcup J_n$ is easily seen to have upper asymptotic density zero, that is $\bar{d}(\bigcup J_n) = 0$ (use (18) and the bound given by (21)). Then we have

$$\bar{d}(x, y) = \bar{d}(\{j \in \omega : x_j \neq y_j\}) = \bar{d}(\bigcup J_n) = 0.$$

Using Lemma 3(a) and the fact that $x \in G_\mu$ we see that $y = \pi(\beta)$ is generic for μ . Hence $\mathcal{C}_3 \subseteq \pi^{-1}(G_\mu)$.

If $\beta \notin \mathcal{C}_3$ then for some strictly increasing sequence of integers (n_k) and some K for each $k \in \mathbb{N}$ we have $\beta(n_k) = K < \infty$. This implies that along the sequence

$(2n_k)$ the frequency of the symbol γ in $y_{[b_{2n_k-1}, b_{2n_k}]}$ is at most $\mu([\gamma])(1 - 1/K) + \varepsilon$ where ε can be made arbitrarily small by choosing k large enough. Thus

$$\liminf_{k \rightarrow \infty} \frac{|\{0 \leq s < b_{2n_k} : y_s = \gamma\}|}{2n_k} < \mu([\gamma]),$$

while using (18), (19) and (20) we get

$$\lim_{k \rightarrow \infty} \frac{|\{0 \leq s < b_{2k+1} : y_s = \gamma\}|}{2n_k + 1} = \mu([\gamma]).$$

This implies that if $\beta \notin \mathcal{C}_3$, then y is an irregular point, $y \in I(X)$. Thus $\pi^{-1}(X \setminus Q(X)) = \pi^{-1}(I(X)) \supseteq \mathbb{N}^{\mathbb{N}} \setminus \mathcal{C}_3$. We conclude $\pi^{-1}(G_\mu) = \pi^{-1}(Q(X)) = \mathcal{C}_3$, and $\pi^{-1}(I(X)) = \mathbb{N}^{\mathbb{N}} \setminus \mathcal{C}_3$. The map π is therefore a reduction map proving that B is Π_3^0 -hard and so G_μ and $Q(X)$ are Π_3^0 -complete and $I(X)$ is Σ_3^0 -complete.

Now suppose $\mu = \delta_0$. Let ν be any ergodic measure on X different from μ and let $x \in G_\nu$. Let $\gamma \neq 0$ be any nonzero symbol such that $\nu([\gamma]) > 0$. Let b_n be an increasing sequence defined as before with μ replaced by ν in (19). Then repeat the definition of auxiliary sets I_n and J_n as above, and define the reduction map $\pi: \mathbb{N}^{\mathbb{N}} \rightarrow X$ by $y = \pi(\beta)$ where

$$y_k = \begin{cases} x_k, & \text{if } k \in \bigcup J_n, \text{ and} \\ 0, & \text{otherwise.} \end{cases}$$

Reasoning as above we see that π is continuous, maps \mathcal{C}_3 into $G_\mu \subseteq Q(X)$, and $\mathbb{N}^{\mathbb{N}} \setminus \mathcal{C}_3$ into $I(X) = X \setminus Q(X) \subseteq X \setminus G_\mu$. This concludes the proof. \square

4. EXAMPLES AND APPLICATIONS

We present here some rather straightforward but noteworthy consequences of Theorem 6. Recall that Ki and Linton [?] showed that in the classical case of r -ary expansions the set of normal numbers is Π_3^0 -complete. We consider several classes of generalized expansions for which our theorem provides a similar result.

Consider first the case of generalized GLS expansions (a generalization of “generalized Lüroth Series”). These include (generalized) Lüroth series expansions, which in turn include r -ary expansions, as well as expansions generated by the tent map. Note that for these applications we can also use Theorem 10 in place of Theorem 6.

4.1. Some generalities. Let $\mathcal{I} = \{I_n = [\ell_n, r_n] \subseteq [0, 1] : n \in \mathcal{D}\}$ be a family of pairwise disjoint intervals indexed by an at most countable set $\mathcal{D} \subseteq \omega$. We call \mathcal{D} the *set of digits* of \mathcal{I} . We assume that \mathcal{I} is a partition of $[0, 1]$ modulo sets of zero Lebesgue measure, that is, we assume $\sum_{n \in \mathcal{D}} (r_n - \ell_n) = 1$. We also set $I_\infty = [0, 1] \setminus \bigcup_{n \in \mathcal{D}} I_n$. Note that $1 \in I_\infty$ and I_∞ may be uncountable. We also define the *address map* $A_{\mathcal{I}}: [0, 1] \rightarrow \mathcal{D} \cup \{\infty\}$ associated with \mathcal{I} by $A_{\mathcal{I}}(x) = k$ if and only if $x \in I_k$, where $k \in \mathcal{D} \cup \{\infty\}$. Given *any* (not necessarily continuous) map $T: [0, 1] \rightarrow [0, 1]$ such that $T|_{\text{int } I_n}$ is continuous and strictly monotone for each $n \in \mathcal{D}$, we define the *itinerary* $\iota(x)$ of $x \in [0, 1]$ with respect to T and \mathcal{I} by $\iota(x) = a_1 a_2 \dots \in (\mathcal{D} \cup \{\infty\})^{\mathbb{N}}$, where $a_n = A_{\mathcal{I}}(T^{n-1}(x))$ for $n \geq 1$. Note that T must be Borel measurable. We say that a Borel probability measure μ on $[0, 1]$ is *T-invariant* if $\mu(B) = \mu(T^{-1}(B))$ for every Borel set $B \subseteq [0, 1]$. A sequence $(x_n)_{n \geq 0} \subseteq [0, 1]$ is *uniformly distributed with respect to μ* if

$$\lim_{N \rightarrow \infty} \frac{1}{N} |\{0 \leq n < N : x_n \in I\}| = \mu(I)$$

for every interval $I \subseteq [0, 1]$ with $\mu(\partial I) = 0$. We say that a point $x \in [0, 1]$ *generates* μ if the sequence $(T^n(x))_{n \geq 0}$ is uniformly distributed with respect to μ .

4.2. Generalized GLS expansions. For more details we refer the reader to the book [?]. Let $\mathcal{I} = \{[\ell_n, r_n) : n \in \mathcal{D}\}$ be a family of intervals as above and fix a function $\epsilon : \mathcal{D} \rightarrow \{0, 1\}$. A *generalized GLS expansion* of $x \in [0, 1]$ determined by (\mathcal{I}, ϵ) is an element $a_1 a_2 \dots$ of $\mathcal{D}^{\mathbb{N}}$ such that

$$(22) \quad x := \frac{h(a_1) + \epsilon(a_1)}{s(a_1)} + \sum_{n=2}^{\infty} (-1)^{\epsilon(a_1) + \dots + \epsilon(a_{n-1})} \frac{h(a_n) + \epsilon(a_n)}{s(a_1)s(a_2) \dots s(a_n)},$$

where $s(n) = 1/(r_n - \ell_n)$ and $h(n) = \ell_n/(r_n - \ell_n)$ for $n \in \mathcal{D}$. Note that for each sequence $a_1 a_2 \dots \in \mathcal{D}^{\mathbb{N}}$ the formula (22) defines a point $x \in [0, 1]$. We write $\psi_{\mathcal{I}, \epsilon}$ for the resulting map from $\mathcal{D}^{\mathbb{N}}$ into $[0, 1]$. Note that $\psi_{\mathcal{I}, \epsilon}$ is continuous, but not necessarily onto. Consider the map $T_{\mathcal{I}, \epsilon}$ such that $T_{\mathcal{I}, \epsilon}(x) = 0$ for each $x \in I_{\infty}$ and on each interval I_n we have that $T_{\mathcal{I}, \epsilon}|_{I_n}$ is a linear function with positive slope from I_n onto $[0, 1]$ if $\epsilon(n) = 0$, and if $\epsilon(n) = 1$ we use the linear map from I_n onto $(0, 1]$ with negative slope. This defines a map $T_{\mathcal{I}, \epsilon} : [0, 1] \rightarrow [0, 1]$. Let I_{∞}^* be the set of all points $x \in [0, 1]$ such that the $T_{\mathcal{I}, \epsilon}$ -orbit of x visits I_{∞} at some iterate, that is $T_{\mathcal{I}, \epsilon}^n(x) \in I_{\infty}$ for some $n \geq 0$. The itinerary map $\iota_{\mathcal{I}, \epsilon}$ determines an (\mathcal{I}, ϵ) -GLS expansion for each $x \in [0, 1] \setminus I_{\infty}^*$. The resulting sequences are called *proper* (\mathcal{I}, ϵ) -GLS expansions and are dense in $\mathcal{D}^{\mathbb{N}}$.

For each x in the set

$$\Omega_{\mathcal{I}, \epsilon} = \bigcap_{k=0}^{\infty} \bigcup_{n \in \mathcal{D}} T_{\mathcal{I}, \epsilon}^{-k}(I_n \cap (0, 1)) = [0, 1] \setminus \bigcup_{k \geq 0} T_{\mathcal{I}, \epsilon}^{-k}(\{0\})$$

the itinerary $\iota_{\mathcal{I}, \epsilon}$ is continuous and gives us the unique (\mathcal{I}, ϵ) -GLS expansion of x . Note that $T_{\mathcal{I}, \epsilon}^{-1}(\{0\}) = [0, 1] \setminus \bigcup_{n \in \mathcal{D}} \text{int } I_n$ is a closed nowhere dense set, hence $\Omega_{\mathcal{I}, \epsilon}$ is a dense G_{δ} set. Furthermore, the function $\iota_{\mathcal{I}, \epsilon}$ is a homeomorphism of $\Omega_{\mathcal{I}, \epsilon}$ onto the set $\iota_{\mathcal{I}, \epsilon}(\Omega_{\mathcal{I}, \epsilon})$ with the inverse given by $\psi_{\mathcal{I}, \epsilon}$ restricted to $\iota_{\mathcal{I}, \epsilon}(\Omega_{\mathcal{I}, \epsilon})$. We also have $\psi_{\mathcal{I}, \epsilon}|_{\iota_{\mathcal{I}, \epsilon}(\Omega_{\mathcal{I}, \epsilon})} \circ \sigma = T_{\mathcal{I}, \epsilon} \circ \psi_{\mathcal{I}, \epsilon}|_{\iota_{\mathcal{I}, \epsilon}(\Omega_{\mathcal{I}, \epsilon})}$. The *fundamental interval* $\Delta(i_1, \dots, i_k)$, where $i_1, \dots, i_k \in \mathcal{D} \cup \{\infty\}$ is the set

$$\{x \in [0, 1] : \iota_{\mathcal{I}, \epsilon}(x) \in [i_1 \dots i_k] \subseteq (\mathcal{D} \cup \{\infty\})^{\mathbb{N}}\}.$$

Fix $i_1, \dots, i_k \in \mathcal{D}$. Take $x \in \Delta(i_1, \dots, i_k)$. Writing p_k/q_k for the partial sum of the (\mathcal{I}, ϵ) -GLS expansion for x given by (22) (summing up to $n = k$) and setting $\epsilon_j = \epsilon(A_{\mathcal{I}}(T_{\mathcal{I}, \epsilon}^{j-1}(x)))$ for $j = 1, \dots, k$ we see that

$$x = \frac{p_k}{q_k} + (-1)^{\epsilon_1 + \dots + \epsilon_k} \frac{T_{\mathcal{I}, \epsilon}^k(x)}{s(i_1) \cdot \dots \cdot s(i_k)}.$$

Since $T_{\mathcal{I}, \epsilon}^k(x)$ takes any value in $[0, 1]$ if $\epsilon_k = 0$ and in $(0, 1]$ if $\epsilon_k = 1$ we have

$$\Delta(i_1, \dots, i_k) = \begin{cases} [d_k, d_k + t_k), & \text{if } \epsilon_k = 0, \text{ and} \\ [d_k - t_k, d_k), & \text{otherwise,} \end{cases}$$

where $d_k = p_k/q_k$, and $t_k = 1/(s(i_1) \cdot \dots \cdot s(i_k))$.

Theorem 11. *Let $T_{\mathcal{I}, \epsilon}$ be the generalized GLS expansion map associated with the pair (\mathcal{I}, ϵ) . If μ is a $T_{\mathcal{I}, \epsilon}$ -invariant Borel probability measure with $\mu(\{0\}) = 0$, then the set of $x \in [0, 1]$ which generate μ is Π_3^0 -complete. Furthermore, the set of irregular points for $T_{\mathcal{I}, \epsilon}$ is Σ_3^0 -complete.*

Proof. First note that $\mathcal{D}^{\mathbb{N}}$ satisfies the assumptions of Theorem 6. Let μ be a $T_{\mathcal{I},\epsilon}$ -invariant Borel probability measure on $[0, 1]$ such that $\mu(\{0\}) = 0$. It follows that $\mu(\Omega_{\mathcal{I},\epsilon}) = 1$. Furthermore, no point in $[0, 1] \setminus \Omega_{\mathcal{I},\epsilon}$ can generate μ , as all these points are eventually mapped to 0 by $T_{\mathcal{I},\epsilon}$. Then we can define $\nu = \iota_{\mathcal{I},\epsilon}^*(\mu)$ and ν is a shift-invariant measure concentrated on $\iota_{\mathcal{I},\epsilon}(\Omega_{\mathcal{I},\epsilon}) \subseteq \mathcal{D}^{\mathbb{N}}$. Since ν is shift-invariant and its support is contained in $\mathcal{D}^{\mathbb{N}}$, which has the specification property (c.f. [?, ?]), the set of ν -generic points G_ν is nonempty and uncountable. On the other hand $\mathcal{D}^{\mathbb{N}} \setminus \iota_{\mathcal{I},\epsilon}(\Omega_{\mathcal{I},\epsilon})$ is at most countable, so $G_\nu \cap \iota_{\mathcal{I},\epsilon}(\Omega_{\mathcal{I},\epsilon}) \neq \emptyset$. For each $z \in G_\nu \cap \iota_{\mathcal{I},\epsilon}(\Omega_{\mathcal{I},\epsilon})$ we have that the σ -orbit of z visits a cylinder $[a_1 \dots a_k]$ with limiting frequency $\nu([a_1 \dots a_k])$ for every $a_1, \dots, a_k \in \mathcal{D}$. Since $z \in \iota_{\mathcal{I},\epsilon}(\Omega_{\mathcal{I},\epsilon})$ and $\psi_{\mathcal{I},\epsilon} \circ \sigma = T_{\mathcal{I},\epsilon} \circ \psi_{\mathcal{I},\epsilon}$ on that set, we have that $\sigma^n(z) \in [a_1 \dots a_k]$ if and only if $T_{\mathcal{I},\epsilon}^n(\psi_{\mathcal{I},\epsilon}(z)) \in \Delta(a_1, \dots, a_k)$. It follows that $\psi_{\mathcal{I},\epsilon}(z)$ visits $\Delta(a_1, \dots, a_k)$ with frequency $\nu([a_1 \dots a_k]) = \iota_{\mathcal{I},\epsilon}^*(\mu)([a_1 \dots a_k]) = \mu(\iota_{\mathcal{I},\epsilon}^{-1}([a_1 \dots a_k]) = \mu(\Delta(a_1, \dots, a_k))$. Furthermore, the boundary points of every basic interval $\Delta(a_1, \dots, a_k)$ are eventually mapped to 0, therefore $\mu(\partial\Delta(a_1, \dots, a_k)) = 0$. Note also that, for each interval J in $[0, 1]$ and $\delta > 0$ we can find a countable family \mathcal{J} of disjoint basic intervals contained in J such that $\mu(J \setminus \bigcup \mathcal{J}) < \delta$. It follows that $\psi_{\mathcal{I},\epsilon}(z)$ generates μ .

Using that $\psi_{\mathcal{I},\epsilon}$ is continuous on $\mathcal{D}^{\mathbb{N}}$ and that $\psi_{\mathcal{I},\epsilon}(z)$ generates μ if and only if $z \in G_\nu \cap \iota_{\mathcal{I},\epsilon}(\Omega_{\mathcal{I},\epsilon})$ we see that to finish the proof we need to show that $G_\nu \cap \iota_{\mathcal{I},\epsilon}(\Omega_{\mathcal{I},\epsilon})$ is Π_3^0 -complete. But this is obvious since G_ν is Π_3^0 -complete by Theorem 6 and $G_\nu \setminus \iota_{\mathcal{I},\epsilon}(\Omega_{\mathcal{I},\epsilon})$ is contained in the set of improper expansions, so it is at most countable.

Now consider any point x which is irregular for $T_{\mathcal{I},\epsilon}$, equivalently, with irregular (\mathcal{I}, ϵ) -GLS expansion. Clearly, $x \in \Omega_{\mathcal{I},\epsilon}$, hence the visits of the $T_{\mathcal{I},\epsilon}$ -orbit of x to some basic interval $\Delta(a_1, \dots, a_k)$, where $a_1, \dots, a_k \in \mathcal{D}$, doesn't have a limiting frequency. It follows that $z = \iota_{\mathcal{I},\epsilon}(x) \in I(\mathcal{D}^{\mathbb{N}})$. By Theorem 6, the irregular set $I(\mathcal{D}^{\mathbb{N}})$ of the full shift on $\mathcal{D}^{\mathbb{N}}$ is Σ_3^0 -complete. Therefore the set of all x irregular for $T_{\mathcal{I},\epsilon}$ equals $\psi_{\mathcal{I},\epsilon}(I(\mathcal{D}^{\mathbb{N}}) \cap \iota_{\mathcal{I},\epsilon}(\Omega_{\mathcal{I},\epsilon}))$. Reasoning as above we see that the latter must be a Σ_3^0 -complete set, which ends the proof. \square

The Lebesgue measure λ on $[0, 1]$ is easily seen to be an invariant ergodic measure for $T_{\mathcal{I},\epsilon}$ (see [?, Chapter 3]). A real $x \in [0, 1]$ is normal for the (\mathcal{I}, ϵ) -GLS expansion if x generates λ .

Corollary 12. *For any (\mathcal{I}, ϵ) -GLS expansion, the set of numbers which are normal for this expansion is Π_3^0 -complete.*

The generalized GLS expansions of Corollary 12 include several types of expansions as we record in the following corollary.

Corollary 13. *Each of the following sets is a Π_3^0 -complete subset of $[0, 1]$: numbers normal for the Lüroth series expansions (see [?]), normal for Q_∞ expansions (see [?]), α -Lüroth expansions (see [?]), and numbers normal for r -ary expansions.*

4.3. β -expansions. Our next application concerns β -expansions. Fix a real number $\beta > 1$. Set $\mathcal{D}_\beta = \{0, 1, \dots, \lfloor \beta \rfloor\}$. For $n \in \mathcal{D}_\beta$ let

$$I_n = \begin{cases} [n/\beta, (n+1)/\beta), & \text{if } 0 \leq n < \lfloor \beta \rfloor, \text{ and} \\ [\lfloor \beta \rfloor/\beta, 1), & \text{otherwise.} \end{cases}$$

Define $\mathcal{I}_\beta = \{I_n : n \in \mathcal{D}_\beta\}$ and $T_\beta(x) = \beta x \bmod 1$ for $x \in [0, 1]$. A β -expansion of $x \in [0, 1]$ is a sequence $d_1 d_2 \dots \in \mathcal{D}_\beta^{\mathbb{N}}$ so that $x = \sum_{i=1}^{\infty} \frac{d_i}{\beta^i}$. For each $x \in [0, 1]$ the

itinerary of x with respect to T_β and \mathcal{I}_β , denoted by $\iota_\beta(x) = d_1 d_2 \dots$ and given by the formula $d_i = \lfloor \beta T_\beta^{i-1}(x) \rfloor$ for $i \geq 1$, defines a sequence $\vec{d} = d_1 d_2 \dots \in \mathcal{D}_\beta^\mathbb{N}$ which is a β -expansion of x . We use the same formula to define the β -expansion of 1, denoted by 1_β . We let $\vec{e} = e_1 e_2 \dots = 1_\beta$ if 1_β does not end in a tail of 0's, and if $1_\beta = d_1 \dots d_k 00 \dots$, where $d_k \neq 0$, then we let $\vec{e} = e_1 e_2 \dots$ be the periodic sequence $(d_1 \dots d_{k-1} d_k - 1)^\infty$. We say a sequence of digits $\vec{d} = d_1 d_2 \dots \in \mathcal{D}_\beta^\mathbb{N}$ is a *proper β -expansion* if there exists $x \in [0, 1)$ such that $\vec{d} = \iota_\beta(x)$. A point $x \in (0, 1)$ has a unique β -expansion \vec{d} given by $\vec{d} = \iota_\beta(x)$ if and only if $T_\beta^i(x) \neq 0$ for each $i \in \mathbb{N}$. If $x \in (0, 1)$ and $T_\beta^i(x) = 0$ for some $i \in \mathbb{N}$, then x has exactly two β -expansions: one proper, and the other we call *improper*. Clearly, the set of improper β -expansions is countable.

A sequence $d_1 d_2 \dots \in \mathcal{D}_\beta^\mathbb{N}$ is *admissible* if it is a β -expansion of some $x \in [0, 1]$. We recall the following well-known fact ([?]): a sequence $d_1 d_2 \dots \in \mathcal{D}_\beta^\mathbb{N}$ is a proper β -expansion if and only if for all $i \in \mathbb{N}$ we have that $d_i d_{i+1} \dots <_{\text{lex}} e_1 e_2 \dots$, where $<_{\text{lex}}$ denotes the strict lexicographic ordering on $\mathcal{D}_\beta^\mathbb{N}$. We note that the sequence \vec{e} itself also has the property that for any shift $\sigma^k(\vec{e}) = e_k e_{k+1} \dots$ we have $\sigma^k(\vec{e}) \leq_{\text{lex}} \vec{e}$. Observe that if \vec{d} is an admissible sequence and \vec{d}' is obtained from \vec{d} by lowering certain digits, then \vec{d}' is also admissible. The set of proper β -expansions $Y_\beta := \iota_\beta([0, 1)) \subseteq \mathcal{D}_\beta^\mathbb{N}$ is shift-invariant but not closed in $\mathcal{D}_\beta^\mathbb{N}$. Let X_β be the closure of Y_β in $\mathcal{D}_\beta^\mathbb{N}$, so X_β is a subshift of $\mathcal{D}_\beta^\mathbb{N}$ known as a *β -shift*. Every β -shift is hereditary. We can characterise X_β as the set of admissible sequences, or equivalently, those sequences $d_1 d_2 \dots \in \mathcal{D}_\beta^\mathbb{N}$ such that for all $i \in \mathbb{N}$ we have that $d_i d_{i+1} \dots \leq_{\text{lex}} \vec{e}$. From this it follows that the set of improper β -expansions $X_\beta \setminus Y_\beta$ is countable and Y_β is a dense G_δ subset of X_β . There is a continuous map $\psi_\beta: X_\beta \rightarrow [0, 1]$ given by $\psi_\beta(d_1 d_2 \dots) = \sum_{i=1}^\infty \frac{d_i}{\beta^i}$. The restriction of the map ψ_β to Y_β is a bijection onto $[0, 1]$, but its inverse $\iota_\beta = (\psi_\beta|_{Y_\beta})^{-1}$ is not continuous on $[0, 1]$. But ι_β is continuous on a subset Ω_β of $[0, 1]$ defined by

$$\Omega_\beta = [0, 1] \setminus \bigcup_{k \geq 0} T_\beta^{-k}(\{0\}).$$

Note that every point in Ω_β has a unique β -expansion, $0 \notin \Omega_\beta$, but 0 has a unique β -expansion, and the only other point in $[0, 1]$ which may have a unique β -expansion and stay outside of Ω_β is 1. Let Z_β be the set of unique β -expansions of points in $(0, 1)$. We have $Z_\beta = \iota_\beta(\Omega_\beta) \subseteq Y_\beta$, more precisely

$$Z_\beta = X_\beta \setminus \{ \{d_1 d_2 \dots \in \mathcal{D}_\beta^\mathbb{N} : (\exists i \geq 1) d_i d_{i+1} \dots = 0^\infty \text{ or } d_i d_{i+1} \dots = \vec{e}\} \}.$$

Thus ι_β restricted to Ω_β is the inverse of $\psi_\beta|_{Z_\beta}$. The admissible sequences can also be described as follows. Let G be the labelled directed graph (with *loops*, that is edges whose initial and terminal vertices are the same) on the vertex set ω defined as follows. Each vertex $i \in \omega$ is the initial vertex of the edge leading to the vertex $i + 1$, and labelled with e_{i+1} . If $e_{i+1} > 0$, then we add e_i many edges from i to 0, and label them with $0, 1, \dots, e_i - 1$. The elements of X_β are obtained by taking an infinite path starting at 0 and reading off the sequence of labels of the edges used to construct the path. The proper β -expansions (elements of Y_β) are exactly the infinite sequences of labels of paths obtained by starting at the vertex 0 and returning to 0 infinitely many times. In particular, $\mathcal{L}(X_\beta)$ corresponds to

the labels of finite paths through G starting at 0. Note that as $e_1 = \lfloor \beta \rfloor$, there are $e_1 > 0$ edges from 0 to 0 (and these are the only loops in the graph G).

Lemma 14. *For every $\beta > 1$ the β -shift X_β has the right feeble specification property.*

Proof. Let $\mathcal{G} \subseteq \mathcal{L}(X_\beta)$ be the set of $w \in \mathcal{D}_\beta^{<\omega}$ corresponding to closed paths in the graph G which start and end at the vertex 0. Clearly if $u \in \mathcal{G}$ and $v \in \mathcal{L}(X_\beta)$ then $uv \in \mathcal{L}(X_\beta)$ (since v corresponds to a label of a path starting at 0). We claim that there is a single symbol in v so that if we change it to 0, then for the resulting word v' we have that $uv' \in \mathcal{G}$, so uv' is a label of a closed path based at the vertex 0. To see this, let $v \in \mathcal{L}(X_\beta)$ with $|v| = k$, and let g_1, \dots, g_k be edges of G whose labels give v . As remarked above we may assume that g_1 starts at the vertex 0. If $v = 0^k$ then there is nothing to prove: we follow the closed path defining u and then use k times the loop based at 0. Otherwise, let $1 \leq i \leq k$ be largest index so that g_i is an edge labelled by a nonzero symbol. It follows that $v = v_1 \dots v_i 0^{k-i}$, where $v_i \neq 0$. Let j be the initial vertex of g_i . Since $v_i \neq 0$ there is an edge from j to 0 labelled with 0. That is, reading off the i th symbol v_i of v , we have the option of returning to 0. Note that $v_j = 0$ for $j > i$. Let v' equal v except we set $v'_i = 0$. Then $uv' \in \mathcal{G}$ as it corresponds to the path through G which starts at 0, returns to 0 at step i , and then loops at 0 until the end of the word. Since we made only one change in v to obtain v' it is easy to see that $d_H(v, v')$ can be arbitrarily small if v is long enough. \square

Theorem 15. *If μ is a T_β -invariant Borel probability measure, then the set of $x \in [0, 1]$ which generate μ is Π_3^0 -complete and the set of points with irregular β -expansion, denoted $I(T_\beta)$, is Σ_3^0 -complete.*

Proof. Let μ be a T_β -invariant measure on $[0, 1]$. Note that $\mu([0, 1)) = 1$, because 1 is never a periodic point for T_β . If $\mu(\{0\}) = 0$, then reasoning as in the proof of Theorem 11 we see that $\mu(\Omega_\beta) = 1$, and $\mu = \psi_\beta^*(\nu)$ for some shift-invariant Borel probability measure supported on $Z_\beta \subseteq Y_\beta$. Note also that $\psi_\beta(0^\infty) = 0$, thus the T_β -invariant Borel probability measure concentrated at 0 is the image through ψ_β^* of the shift-invariant Borel probability measure concentrated at $0^\infty \in Y_\beta$. It follows that every T_β -invariant Borel probability measure on $[0, 1]$ is the image through ψ_β^* of a shift-invariant measure on Y_β .

Fix a T_β -invariant measure μ on $[0, 1]$. By the above there is a shift-invariant measure ν on Y_β such $\mu = \psi_\beta^*(\nu)$. Observe that $\nu(Y_\beta) = 1$ implies that $\nu(\{\vec{e}\}) = 0$. Given $a_1 \dots a_k \in \mathcal{L}(X_\beta)$, define the basic interval

$$\Delta(a_1, \dots, a_k) = \{x \in [0, 1] : \iota_\beta(x) \in [a_1 \dots a_k]\}.$$

Since $\mu = \psi_\beta^*(\nu)$, we have

$$\mu(\Delta(a_1, \dots, a_k)) = \nu(\psi_\beta^{-1}(\Delta(a_1, \dots, a_k)))$$

and

$$\psi_\beta^{-1}(\Delta(a_1, \dots, a_k)) = \left([a_1 \dots a_k] \cap Y_\beta\right) \cup \left([a_1 \dots a_k] \cap \bigcup_{n \geq 0} \sigma^{-n}(\{\vec{e}\})\right).$$

It follows that

$$(23) \quad \mu(\Delta(a_1, \dots, a_k)) = \nu([a_1 \dots a_k] \cap Y_\beta) = \nu([a_1 \dots a_k]).$$

Note also that for every T_β -invariant Borel probability measure μ and a basic interval $\Delta(a_1, \dots, a_k)$ we have that

$$\partial\Delta(a_1, \dots, a_k) \subseteq \bigcup_{n \geq 1} T_\beta^{-n}(\{0\}) \cup \{1\}.$$

(Remember that 0 is an interior point of any interval $[0, r)$, where $r > 0$.) Since basic intervals generate the Borel sigma algebra, we see that a point $x \in [0, 1)$ generates μ if and only if for every $a_1 \dots a_k \in \mathcal{L}(X_\beta)$ the T_β -orbit of x visits the basic interval $\Delta(a_1, \dots, a_k)$ with frequency $\mu(\Delta(a_1, \dots, a_k))$. Observe that if $a_1 \dots a_k \in \mathcal{L}(X_\beta)$, $z \in Y_\beta$, and $\sigma^n(z) \in [a_1 \dots a_k] \subseteq X_\beta$, then $T_\beta^n(\psi_\beta(z))$ belongs to the basic interval $\Delta(a_1, \dots, a_k)$, since we have $\psi_\beta \circ \sigma = T_\beta \circ \psi_\beta$ on Y_β . In particular, if $z \in Y_\beta$ is generic for ν , then using (23) we see that $\psi_\beta(z)$ visits a basic interval $\Delta(a_1, \dots, a_k)$ with frequency $\nu([a_1 \dots a_k])$, so $\psi_\beta(z) \in [0, 1)$ generates μ . Conversely, if x generates μ , then the T_β -orbit of x visits each basic interval $\Delta(a_1, \dots, a_k)$ with frequency $\mu(\Delta(a_1, \dots, a_k))$, which means that the orbit of $\iota_\beta(x)$ under σ visits the cylinder set $[a_1 \dots a_k]$ with the same frequency, so $\iota_\beta(x)$ is generic for ν on Y_β . It follows that $\psi_\beta(G_\nu \cap Y_\beta)$ is the set of points in $[0, 1)$ that generate μ .

By Lemma 14 the subshift X_β has the right feeble specification property. It is also known that the set of shift-invariant Borel probability measures supported on X_β is uncountable. Thus, X_β satisfies the assumptions of Theorem 6, and we conclude that for each shift-invariant Borel probability measure ν supported on X_β the set $G_\nu \subseteq X_\beta$ of generic points for ν is Π_3^0 -complete. Since $G_\nu \setminus (G_\nu \cap Y_\beta)$ is at most countable, we see that $G_\nu \cap Y_\beta$ is also a Π_3^0 -complete set, and ψ_β reduces it to the set of points in $[0, 1)$ that generate μ . Thus the latter set is also Π_3^0 -complete.

Let $I(X_\beta)$ be the set of irregular points for X_β . Using Theorem 6 again, we see that $I(X_\beta)$ is a Σ_3^0 -complete subset of X_β . Then $I(X_\beta) = (I(X_\beta) \cap X_\beta \setminus Z_\beta) \cup (I(X_\beta) \cap Z_\beta)$ is a disjoint union and $(I(X_\beta) \cap X_\beta \setminus Z_\beta)$ is at most countable. Hence $(I(X_\beta) \cap X_\beta \setminus Z_\beta)$ is a Π_3^0 -set. If $I(X_\beta) \cap Z_\beta$ were also a Π_3^0 -set, then $I(X_\beta)$ would not be Σ_3^0 -complete, which is absurd. Thus $I(X_\beta) \cap Z_\beta$ is a Σ_3^0 -complete set. Reasoning as above we also obtain that $\psi_\beta(I(X_\beta) \cap Z_\beta) = I(T_\beta) \setminus \{T_\beta^n(1) : n \geq 0\}$, which implies that $I(T_\beta)$ is a Σ_3^0 -complete set. \square

For each $\beta > 1$ there is a Borel probability measure on $[0, 1)$ which is invariant for the transformation T_β , which is known as *Parry measure*. It is characterised as the unique ergodic T_β -invariant that is equivalent to Lebesgue measure on $[0, 1)$. We let μ_β denote the Parry measure on $[0, 1)$. A real number $x \in [0, 1)$ is *normal with respect to the β -expansion* if x generates μ_β .

Corollary 16. *For each $\beta > 1$ the set of $x \in [0, 1)$ which are normal with respect to the β -expansion is a Π_3^0 -complete set.*

4.4. Continued fraction expansions. Our next application of Theorem 6 is to the regular continued fraction expansion. Let $T: [0, 1] \setminus \mathbb{Q} \rightarrow [0, 1] \setminus \mathbb{Q}$ be the continued fraction map given by $T(x) = \frac{1}{x} - \lfloor \frac{1}{x} \rfloor$. Let μ be the Gauss measure on $[0, 1]$, which is defined by $\mu(A) = \frac{1}{\ln(2)} \int_0^1 \frac{\chi_A(x)}{1+x} dx$. The Gauss measure is a T -invariant ergodic measure equivalent to the Lebesgue measure. If we let $d(x) = \lfloor \frac{1}{x} \rfloor$, then the regular continued fraction expansion of x is given by $d_1 d_2 \dots \in \mathbb{N}^\mathbb{N}$, where $d_i = d(T^{i-1}(x))$ for $i \in \mathbb{N}$. This expansion gives a homeomorphism ι between

$[0, 1] \setminus \mathbb{Q}$ and $\mathbb{N}^{\mathbb{N}}$ such that $\iota \circ T = \sigma \circ \iota$, where σ is the shift map on $\mathbb{N}^{\mathbb{N}}$. Every T -invariant Borel probability measure μ on $[0, 1] \setminus \mathbb{Q}$ corresponds to a unique shift-invariant Borel probability measure $\nu = \iota^*(\mu)$ on $\mathbb{N}^{\mathbb{N}}$. Since the full shift $\mathbb{N}^{\mathbb{N}}$ satisfies the assumptions of Theorem 6 we see that the set of sequences generic for ν , denoted by G_ν , is a Π_3^0 -complete subset of $\mathbb{N}^{\mathbb{N}}$. It follows that the set of points that generate μ given by $\iota^{-1}(G_\nu)$ is also Π_3^0 -complete. The same reasoning shows that the set of continued fraction irregular points is Σ_3^0 -complete.

Theorem 17. *If μ is a Borel probability measure on $[0, 1] \setminus \mathbb{Q}$ which is invariant for the continued fraction map, then the set of points in $[0, 1] \setminus \mathbb{Q}$ that generate μ is a Π_3^0 -complete set and the set of points with irregular continued fraction expansion is Σ_3^0 -complete.*

Corollary 18. *The set of $x \in [0, 1] \setminus \mathbb{Q}$ which are continued fraction normal is a Π_3^0 -complete set.*

4.5. Sharkovsky-Sivak problem and the tent map. Our last application considers the tent map $T: [0, 1] \rightarrow [0, 1]$ given by

$$T(x) = \begin{cases} 2x, & \text{if } 0 \leq x \leq \frac{1}{2}, \text{ and} \\ 2 - 2x, & \text{if } \frac{1}{2} < x \leq 1. \end{cases}$$

Taking $\mathcal{I} = \{I_0 = [0, 1/2), I_1 = [1/2, 1)\}$ and the function $\epsilon: \{0, 1\} \rightarrow \{0, 1\}$ such that $\epsilon(0) = 0$ and $\epsilon(1) = 1$ we can easily see that $T_{\mathcal{I}, \epsilon}$ is the tent map and the (\mathcal{I}, ϵ) -GLS expansion map coincides with the tent map. Moreover, it is well-known and easy to see, either directly or following the reasoning presented above for the general GLS expansions, that the tent map is a factor of the full shift system $\{0, 1\}^{\mathbb{N}}$ under a factor map $\psi_{\mathcal{I}, \epsilon}$ which is onto and one-to-one except at the countable set $\bigcup_{n \geq 0} T^{-n}(1/2)$, where $\psi_{\mathcal{I}, \epsilon}$ is two-to-one (see also [?], Example E in 6.3.5, taking into account Proposition 6.3.4 (2) therein). As a corollary we obtain the following result.

Corollary 19. *If μ is a Borel probability measure invariant for the tent map T , then the set of points that generate μ (also known as the statistical basin for μ) is a Π_3^0 -complete set. The set of irregular points is Σ_3^0 -complete.*

In particular, the statistical basin for the Dirac mass at 0 and the tent map is a Π_3^0 -complete set, which answers [?, Problem 5].

Also as a corollary of Theorem 6 we can answer a question of Sharkovsky and Sivak [?, Problem 3], who asked whether there is a continuous map $f: [0, 1] \rightarrow [0, 1]$ which has an invariant Borel probability measure μ such that the set of generic points in Π_3^0 -complete.

5. CONCLUDING REMARKS

Note that there are numeration systems for which our approach does not work. For example, the Cantor series expansions are obtained through nonautonomous dynamical systems and thus require a separate analysis. The most up to date and general results on normal numbers in this context are found in [?, ?, ?], respectively. In [?] descriptive complexity results similar to the ones in the present paper are obtained for Cantor series expansions. With the results presented here, this shows that Π_3^0 -completeness is another universal property that holds for all known examples of normal numbers.

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(D. Airey) DEPARTMENT OF MATHEMATICS, UNIVERSITY OF TEXAS AT AUSTIN, 2515 SPEEDWAY, AUSTIN, TX 78712-1202, USA

Current address, D. Airey: Department of Mathematics, Princeton University, Fine Hall, Washington Road, Princeton, NJ 08544-1000, USA

E-mail address: dairey@math.princeton.edu

(S. Jackson) DEPARTMENT OF MATHEMATICS, UNIVERSITY OF NORTH TEXAS, GENERAL ACADEMICS BUILDING 435, 1155 UNION CIRCLE, #311430, DENTON, TX 76203-5017, USA

E-mail address: stephen.jackson@unt.edu

(D. Kwietniak) FACULTY OF MATHEMATICS AND COMPUTER SCIENCE, JAGIELLONIAN UNIVERSITY IN KRAKOW, UL. ŁOJASIEWICZA 6, 30-348 KRAKÓW, POLAND, AND INSTITUTE OF MATHEMATICS, FEDERAL UNIVERSITY OF RIO DE JANEIRO, CIDADE UNIVERSITARIA - ILHA DO FUNDÃO, RIO DE JANEIRO 21945-909, BRAZIL

E-mail address: dominik.kwietniak@uj.edu.pl

URL: www.im.uj.edu.pl/DominikKwietniak/

(B. Mance) INSTITUTE OF MATHEMATICS OF POLISH ACADEMY OF SCIENCE, ŚNIADECKICH 8, 00-656 WARSAW, POLAND

Current address, B. Mance: Uniwersytet im. Adama Mickiewicza w Poznaniu, Collegium Mathematicum, ul. Umultowska 87, 61-614 Poznań, Poland

E-mail address: william.mance@amu.edu.pl