

FIRST RETURN SYSTEMS FOR SOME CONTINUED FRACTION MAPS

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ABSTRACT. We prove a conjecture of Calta, Kraaikamp and the author: For all $n \geq 3$, each member of their one-parameter family of interval maps, denoted $T_{3,n,\alpha}$, has its ‘first expansive return map’ of natural extension given by the first return map under the geodesic flow to a section of the unit tangent bundle of the hyperbolic surface uniformized by the underlying Fuchsian group $G_{3,n}$.

To achieve the proof, we first prove the corresponding result for analogous one-parameter families related to the Hecke triangle Fuchsian group $G_{2,n}$. A direct comparison per n of the $\alpha = 1$ planar domains allows the Hecke group setting to provide sufficient information to prove the conjecture.

We also give details about the entropy functions for the Hecke triangle Fuchsian group maps, $\alpha \mapsto h(T_{2,n,\alpha})$. Each is continuous on $(0, 1)$, increasing on $(0, 1/2)$, decreasing on $(1/2, 1)$, with a central interval of constancy. We give precise formulas for the end points of the central intervals and also give precise formulas for the maximal entropy per family. For fixed α , the entropy of $T_{2,n,\alpha}$ goes to zero as n tends to infinity.

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Date: 9 November 2025.

2010 Mathematics Subject Classification. 37A25, (11K50, 37A10, 37A44).

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1. INTRODUCTION

1.1. A conjecture proven, and results on entropy functions in the $(2, n)$ setting. The Fuchsian triangle groups with one cusp and two elliptic conjugacy classes are, up to conjugation, of the form $G_{m,n}$ with integers $n \geq m \geq 2$ (where if $m = 2$ then $n > m$), where $G_{m,n}$ is the group generated by

$$(1) \quad A = \begin{pmatrix} 1 & t \\ 0 & 1 \end{pmatrix}, B = \begin{pmatrix} \nu & 1 \\ -1 & 0 \end{pmatrix}, C = \begin{pmatrix} -\mu & 1 \\ -1 & 0 \end{pmatrix},$$

for $\mu = \mu_m = 2 \cos \pi/m$, $\nu = \nu_n = 2 \cos \pi/n$ as well as $t = \mu + \nu$.

For each of these groups and each $\alpha \in [0, 1]$, [CKS] defined the interval map $T_{m,n,\alpha}$ given in (4) below. These maps are studied in [CKS, CKS3, CKS2] with emphasis on the families $T_{3,n,\alpha}$. To an eventually expansive interval map, [CKS3] associate its *first pointwise expansive power map*, see § 2.9 below. They conjectured that the first pointwise expansive power map associated to each $T_{3,n,\alpha}$ is given by the first return map of the geodesic flow to a cross-section of the unit tangent bundle of the hyperbolic surface uniformized by $G_{3,n}$. They proved the conjecture for $n = 3$ and in general reduced the conjecture to showing that the volume of the unit tangent bundle equals product of the Sinai-Kolmogorov measure theoretic entropy (hereafter simply: entropy) $h(T_{3,n,\alpha})$ with $\mu(\Omega_{3,n,\alpha})$, where $\Omega_{3,n,\alpha}$ is a planar domain for a 2-dimensional map naturally associated to $T_{3,n,\alpha}$ and μ is the measure defined in (9) below. We prove the conjecture by showing this equality.

Theorem 1. *For all $n \geq 3$ and for all $\alpha \in (0, 1)$, the product $h(T_{3,n,\alpha})\mu(\Omega_{3,n,\alpha})$ equals the volume of the unit tangent bundle of the hyperbolic orbifold uniformized by the triangle group $G_{3,n}$. Furthermore, the first pointwise expansive power of $T_{3,n,\alpha}$ has its natural extension given by the first return of the geodesic flow to a cross section in the unit tangent bundle of the hyperbolic orbifold uniformized by $G_{3,n}$.*

Note that throughout the remainder of this paper, we will simply write “surface” instead of orbifold.

To reach these results, we prove similar results for the natural analogs of our maps, denoted $T_{2,n,\alpha}$, but related to the Hecke triangle groups $G_{2,n}$. Beyond what is strictly necessary to prove the conjecture, we also have the following. For $n \geq 3$ let R_n be the positive root of $x^2 - (2 - t_{2,n})x - 1$ if n is odd, and for even n let $R_n = 1$.

Theorem 2. *For each $n \geq 3$ the entropy function $\alpha \mapsto h(T_{2,n,\alpha})$ is finite and nonzero on $0 < \alpha < 1$. Furthermore, this function is: increasing on $0 < \alpha < 1 - \frac{R_n}{t_{2,n}}$; is constant on $[1 - \frac{R_n}{t_{2,n}}, 1 + \frac{R_n}{t_{2,n}})$; and, decreases on $1 + \frac{R_n}{t_{2,n}} \leq \alpha < 1$. Its maximal value is*

$$\begin{cases} \frac{\pi^2 (n-2)}{2n \ln \left(\frac{1 + \cos \pi/n}{\sin \pi/n} \right)} & \text{if } 2|n; \\ \frac{\pi^2 (n-2)}{2n \ln \left(\frac{1 + R_n}{2 \sin \frac{\pi}{2n}} \right)} & \text{otherwise.} \end{cases}$$

Finally, for any $0 < \alpha < 1$, the limit as $n \rightarrow \infty$ of $h(T_{2,n,\alpha})$ equals zero.

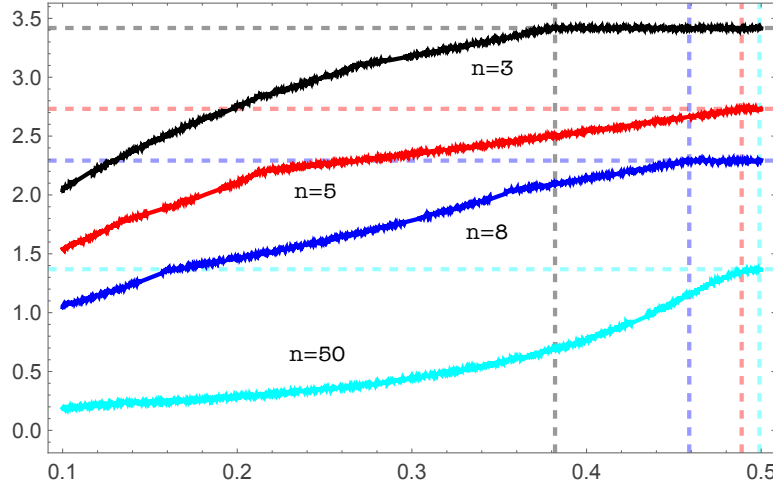


FIGURE 1. Birkhoff sums approximation to entropy functions $\alpha \mapsto h(T_{2,n,\alpha})$, for $n \in \{3, 5, 8, 50\}$ and $0.1 < \alpha \leq 1/2$. Horizontal dashed lines at maximum value as per Theorem 2; vertical dashed lines at $1 - R_n/t_{2,n}$ per n .

1.2. Motivation. For much of the technical vocabulary mentioned here, see § 2 Background.

Although the connection between the geodesic flow on the unit tangent bundle of the modular surface and the regular continued fraction interval map goes back to at least Artin's [Art] in the 1920s, it was only in the 1980s and 1990s that one had proofs [AF, S, Ar] that this interval map was a factor of the dynamical system defined by the first return under the geodesic flow to a section of the unit tangent bundle of the modular surface.

In the ensuing years, various works have shown similar results for other interval maps also given piecewise by the fractional linear action of elements of a Fuchsian group of finite covolume. Among these are [AF2, AF3, N2, MS, KU3, BM, AK, AS]. In these works, consideration of the

Sinai-Kolmogorov measure theoretic entropy (hereafter simply: entropy), $h(T)$, of the interval map T arises naturally.

The behavior of the entropy function for each of a parametrized family of continued fraction-like interval maps was studied by Nakada [N] in the setting of what are now known as the Nakada α -continued fraction maps. These maps, T_α with $\alpha \in [0, 1]$, are all piecewise by elements of the modular group $\mathrm{PSL}_2(\mathbb{Z})$ acting linearly fractionally. Nakada and Natsui [NN] pointed to a phenomenon that now goes under the name of ‘matching intervals’, these are intervals of α such that $T_\alpha^k(\alpha) = T_\alpha^{k'}(\alpha - 1)$ for pairs of positive integers (k, k') ; the function $\alpha \mapsto h(T_\alpha)$ is increasing, decreasing or constant on the interval according to whether $k > k'$, $k < k'$, or $k = k'$. Luzzi and Marmi [LM] strongly suggested that $\alpha \mapsto h(T_\alpha)$ is a continuous function of α . The continuity was proven a few years later by both Carminati and Tiozzo [CT] and [KSS]. Luzzi-Marmi also asked if every T_α is the factor of some cross section to the geodesic flow on the unit tangent bundle of the modular surface. This was proven in [AS], using the techniques of the type sketched below in § 2.6.

After Calta and the author [CS] considered various metric properties of the interval maps $T_{3,n,0}$, Calta, Kraaikamp and this author [CKS] then introduced the $T_{m,n,\alpha}$ with concentration on dynamics of the case of $m = 3$. By the way, the reason we avoided the case of $m = 2$ was so as to not overlap with results already in the literature. Continuity of the entropy functions $\alpha \mapsto h(T_{3,n,\alpha})$ is shown for each $n \geq 3$ in [CKS3], using results of [CKS2].

A naive hope could be that any somewhat natural interval map T which is given piecewise by the fractional linear action of a reasonable set of elements of a Fuchsian group of finite covolume should arise as the factor of the first return system of the geodesic flow on the unit tangent bundle of the hyperbolic surface uniformized by the group. However, already [AS] gave an example of an interval map which, although given piecewise by the fractional linear action of elements of the modular, is not (as discoverable by the techniques of § 2.6) of this type. Still, the maps $T_{3,n,\alpha}$ are quite naturally, one could even say geometrically, defined in terms of elements of $G_{3,n}$. On the other hand, these maps (in general) are only eventually expansive and the methods we use demand true expansiveness. We were thus lead to introduce an acceleration of the maps. The ‘first pointwise expansive power’ of a map as defined in [CKS3] is recalled in § 2.9 below. With this in place, [CKS3] made the conjecture that each $T_{3,n,\alpha}$ has its first pointwise expansive power of natural extension given by the first return by the geodesic flow to a cross-section in the unit tangent bundle of the surface uniformized by $G_{3,n}$.

1.3. Progress in [CKS3] towards proof. A proof in the case of $n = 3$ and all values of α is given in [CKS3]. They also give numerical evidence for the conjecture for larger values of n , but that evidence is questionable already for small n , see [[CKS3], § 15.4]. We quote from that paper: “The integrals which appear are naturally related to dilogarithmic functions and hence are notoriously difficult to evaluate exactly. See [F] for discussion of the appearance of the dilogarithm in a related setting.”

The conjecture (for general n and $0 < \alpha < 1$) is reduced in [CKS3] in two steps. Already to prove the continuity of the entropy functions, for each $n \geq 3$ and $0 < \alpha < 1$ they give a planar extension $\mathcal{T}_{3,n,\alpha} : \Omega_{3,n,\alpha} \rightarrow \Omega_{3,n,\alpha}$, whose action on its first, thus x -, coordinate is given by $T_{3,n,\alpha}$. See § 2.3 for the basics of this type of planar extension. The first reduction step shows that the conjectured result holds for all $n \geq 3$ and all $0 < \alpha < 1$ for which $h(T_{3,n,\alpha})\mu(\Omega_{3,n,\alpha}) = \mathrm{vol}(T^1(G_{3,n} \backslash \mathbb{H}))$, where this latter denotes the volume of the unit tangent bundle of the surface uniformized by $G_{3,n}$, see (13). This reduction step relies on: Rohlin’s entropy formula for entropy, see (6); the fact that the invariant measure for each $T_{3,n,\alpha}$ is the marginal measure on $\Omega_{3,n,\alpha}$; and, what we call ‘Arnoux’s method’, see § 2.6. In brief, this last gives a measure preserving map sending (in the setting discussed) $\Omega_{3,n,\alpha}$ to $T^1(G_{3,n} \backslash \mathbb{H})$ such that the integrand

$\ln |T'(x)|$ in Rohlin's formula is also the length of the geodesic path following the geodesic flow from the image of $(x, y) \in \Omega_{3,n,\alpha}$ to the point corresponding to $\mathcal{T}_{3,n,\alpha}(x, y)$. The second step replaces showing the equality for each n and all $0 < \alpha < 1$ of $h(T_{3,n,\alpha})\mu(\Omega_{3,n,\alpha})$ with the volume by showing the equality to this volume of the integral over (the infinite mass) $\Omega_{3,n,0}$ with respect to μ of the Rohlin integrand (for the map $T_{3,n,0}$). We call such integrals *Rohlin integrals*. That is, they show that the a priori varying values $h(T_{3,n,\alpha})\mu(\Omega_{3,n,\alpha})$ are all equal to the Rohlin integral over $\Omega_{3,n,0}$. See Lemma 3.12 for an analogous result when $m = 3$ is replaced by $m = 2$.

The proof of the conjecture in the case of $m = n = 3$ in [CKS3] proceeds by partitioning $\Omega_{3,3,0}$ — determined in [CS] — into the union of two rectangles, with the value of the Rohlin integral being separately recognizable as equalling the volume of the unit tangent bundle of the modular surface, $\text{vol}(T^1(G_{2,3}\backslash\mathbb{H}))$. Again from the well-known formula (13), arithmetic then gives the result.

1.4. Impetus for this paper. On the occasion of a conference in April 2025 at the E. Schröder Institute in Vienna which was held in part to celebrate the retirement of C. Kraaikamp, the author decided to revisit the conjecture of [CKS3].

1.5. Approach in this paper. Despite the aforementioned earlier hesitation of treating families of interval maps the Hecke group setting, we prove the conjecture by first showing that the analogous conjecture when $m = 2$ holds for any n and $0 < \alpha < 1$. This path was taken because of an observation that each $\Omega_{3,n,1}$ is the union of a rectangle whose Rohlin integral equals $\text{vol}(T^1(G_{2,3}\backslash\mathbb{H}))$ and its complement, with computation suggesting that this second piece would be given by a simple transformation applied to what was sure to be $\Omega_{2,n,1}$. Furthermore, several years ago Giulio Tiozzo encouraged the author to consider the various $T_{2,n,\alpha}$. Recollection of that interaction fully increased our belief that the work of [CKS, CKS3, CKS2] would all go through in the $m = 2$ setting. The key extra ingredient is a result of Nakada [N3], Theorem 2.3 below, which gives a relation of the entropy of each of the Rosen fraction maps to $\text{vol}(T^1(G_{2,3}\backslash\mathbb{H}))$. As [AS2] used without proof, in the notation of this paper, this leads directly to $h(T_{2,n,1/2})\mu(\Omega_{2,n,\alpha}) = \text{vol}(T^1(G_{2,n}\backslash\mathbb{H}))$. We give a proof in § 3.2.

1.6. Entropy maps in the $m = 2$ setting. We are purposely light throughout the paper in our treatment of the case $(m, n) = (2, 3)$. In fact, the $T_{2,3,\alpha}$ are maps within the family of $G_{2,3} = \text{PSL}_2(\mathbb{Z})$ -related ‘ (a, b) -continued fraction maps’ introduced by Katok and Ugarcovici [KU2]. The $T_{2,3,\alpha}$ themselves, and in particular the entropy function $\alpha \mapsto h(T_{2,3,\alpha})$, were studied in [CIT]. Key to their approach is also a matching condition, which can be compared to be the matching condition of Lemma 3.3 (iv). Their approach relies on ‘Farey words’ in two letters to describe the matching intervals. One can expect a direct formula to pass from their two variable labeling of the matching intervals to our (k, v) labels, but we do not pursue that here. We do give Figure 4 to allow the reader to check our planar domain against one figured in [CIT], and an example to support our figure. That all said, it was rereading the work of [CIT] which suggested to us to take the detour to achieve Theorem 2.

We used Birkhoff sum approximations to find the approximate curves for the $\alpha \mapsto h(T_{2,n,\alpha})$ shown in Figure 1. We first saw such ideas in [C], and in the setting of continued fraction-like interval maps in [LM]. See [T] for results of convergence of Birkhoff sum approximations in a very related setting. In fact, the Birkhoff sum approximations led to discovering the error in [BKS] which we point out and correct in § 3.3. Note that Figure 1 suggests that there is a change in concavity of the functions as n goes from small to large values; we have not pursued that matter here.

1.7. Outline. Section 2 is a rather long background section, where we give elementary definitions and reminders through the less elementary matters of Arnoux's method and the setting of [CKS, CKS3, CKS2], as well as state a theorem of Nakada [N3]. Section 3 presents the results in the setting of $m = 2$; the section begins with a statement of results and then a somewhat detailed outline of its subsections. Section 4 gives the shift from $\Omega_{2,n,1}$ to $\Omega_{3,n,1}$ and its implications. Section 5 gives the final lines of the proof of Theorem 1, thus showing that the conjecture from [CKS3] holds.

2. BACKGROUND

We give background, including repeating background from [CKS, CKS3, CKS2].

2.1. The groups, the intervals, and the functions. We consider the following triangle Fuchsian groups. Fix integers $n \geq m \geq 2$ (where if $m = 2$ then $n > m$), and let $\mu = \mu_m = 2 \cos \pi/m$, $\nu = \nu_n = 2 \cos \pi/n$ as well as $t = \mu + \nu$. Thus,

$$t := t_{m,n} = 2 \cos \pi/m + 2 \cos \pi/n.$$

We denote by $G_{m,n}$ the group generated by

$$(2) \quad A = \begin{pmatrix} 1 & t \\ 0 & 1 \end{pmatrix}, B = \begin{pmatrix} \nu & 1 \\ -1 & 0 \end{pmatrix}, C = \begin{pmatrix} -\mu & 1 \\ -1 & 0 \end{pmatrix},$$

and note that $C = AB$. We work projectively, hence B, C are of order n, m respectively while A is of infinite order. That is, $G_{m,n}$ is a Fuchsian triangle group of signature $(0; m, n, \infty)$.

Fix $\alpha \in [0, 1]$ and define

$$(3) \quad \mathbb{I}_{m,n,\alpha} := \mathbb{I}_\alpha = [(\alpha - 1)t, \alpha t].$$

Let

$$(4) \quad T_\alpha = T_{m,n,\alpha} : x \mapsto A^k C^l \cdot x,$$

where as usual, any 2×2 matrix $\begin{pmatrix} a & b \\ c & d \end{pmatrix}$ acts on real numbers by $\begin{pmatrix} a & b \\ c & d \end{pmatrix} \cdot x = \frac{ax+b}{cx+d}$, and

- $l \in \mathbb{N}$ is minimal such that $C^l \cdot x \notin \mathbb{I}$;

•

$$k = -\lfloor (C^l \cdot x)/t + 1 - \alpha \rfloor.$$

We consider T_α as a map on the closed interval taking values in the half-open interval $\mathbb{I}_\alpha$,

$$T_\alpha : [(\alpha - 1)t, \alpha t] \rightarrow \mathbb{I}_\alpha.$$

When $\alpha = 0$ and $(m, n) = (3, n)$ this gives the (unaccelerated) maps treated in [CS], the maps with $(m, n) = (3, n)$ and $0 < \alpha < 1$ are studied in some detail in [CKS, CKS2]. That the $T_{m,n,\alpha}$ lead to continued fraction-like expansions of real numbers is shown in [[CKS], § 1.5].

2.2. Two fundamental formulas for entropy. The Kolmogorov–Sinai measure theoretic entropy, which as stated above we refer to simply as entropy and usually denote in the form $h(T)$ — is an invariant of a dynamic system, which roughly speaking measures its complexity. Rohlin introduced the notion of the natural extension system to aid in the study of entropy and showed that the original system and its natural extension share entropy values.

For a dynamical system $(X, T, \mathcal{B}, \mu)$ of finite measure and a subset $E \subset X$ of positive measure, the *induced transformation* on E is $T_E : E \rightarrow E$ given by $T_E(x) = T^k(x)$ where $k \in \mathbb{N}$ is minimal such that $T^k(x) \in E$. (By the Poincaré Recurrence Theorem, the set of $x \in E$ such that there is some such k has full measure in E , and in fact one defines T_E to be the identity on the complement of this subset.) We set μ_E to be the restriction to E of μ scaled by $1/\mu(E)$. This allows one to define a dynamical system for T_E . The following is a key tool in the study of planar extensions. *Abramov’s Formula* states that the entropy of the induced system on E is the quotient of the entropy of the original system divided by the measure of E , in short

$$(5) \quad h(T_E) = h(T)/\mu(E).$$

One also has *Rohlin’s Entropy Formula* for an interval map $T : \mathbb{I} \rightarrow \mathbb{I}$ which is ergodic with respect to an invariant probability measure ν and such that the derivative T' exists ν -almost everywhere, (see say [DKU]):

$$(6) \quad h(T) = \int_{\mathbb{I}} \ln |T'(x)| d\nu.$$

2.3. Piecewise Möbius maps, cylinders, planar extensions. For

$$(7) \quad M = \begin{pmatrix} a & b \\ c & d \end{pmatrix}$$

in $\mathrm{SL}(2, \mathbb{R})$ and $z \in \mathbb{C}$, one has $M \cdot z = \frac{az+b}{cz+d}$. This fractional linear action allows us to identify M with a Möbius function. We use mainly the restriction of this action to the real numbers. Note that M and $-M$ act identically, thus we may consider that the group $\mathrm{PSL}(2, \mathbb{R})$, the quotient group of $\mathrm{SL}(2, \mathbb{R})$ by $\{\pm I\}$, is acting.

We consider various functions T on a subinterval $\mathbb{I} \subset \mathbb{R}$ with values in $\mathbb{I}$ such that there is a partition $\mathbb{I} = \cup_{\beta \in B} K_\beta$ with $T(x) = M_\beta \cdot x$ for all $x \in K_\beta$, where each M_β is in $\mathrm{SL}_2(\mathbb{R})$ and acts fractionally linearly. (In general, one can also allow M of negative determinants, but this is unnecessary in this paper.) We will assume that each K_β is an interval and is taken as large as possible. We call these K_β the (rank one) *cylinders* for T . Similarly, a cylinder of rank $m > 1$ is the largest interval on which T^m is given by the action of some $M_{\beta_m} \cdots M_{\beta_1}$. We say that the cylinder K_β is *full* if $T(K_\beta) = \mathbb{I}$.

The standard number theoretic planar map associated to $M \in \mathrm{SL}_2(\mathbb{R})$ is

$$\mathcal{T}_M(x, y) := \begin{pmatrix} M \cdot x, RMR^{-1} \cdot y \end{pmatrix} \quad \text{for } x \in \mathbb{R} \setminus \{M^{-1} \cdot \infty\}, y \in \mathbb{R} \setminus \{(RMR^{-1})^{-1} \cdot \infty\},$$

where

$$(8) \quad R = \begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix}.$$

An elementary Jacobian matrix calculation verifies that the measure μ on $\mathbb{R}^2$ given by

$$(9) \quad d\mu = \frac{dx dy}{(1 + xy)^2}$$

is (locally) $\mathcal{T}_M$ -invariant.

For a piecewise Möbius interval map T we then set

$$(10) \quad \mathcal{T}(x, y) = \mathcal{T}_{M_\beta}(x, y) \quad \text{whenever } x \in K_\beta, y \in \mathbb{R} \setminus \{N^{-1} \cdot \infty\}.$$

Suppose that $\Omega \subset \mathbb{R}^2$ projects onto the interval $\mathbb{I}$ and is a domain of bijectivity of $\mathcal{T}$, that is $\mathcal{T}$ is bijective on Ω up to μ -measure zero. Let $\mathcal{B}$ be the Borel algebra of Ω , we then call the system $(\mathcal{T}, \Omega, \mathcal{B}, \mu)$ a *planar extension* for T . We will occasionally abuse this terminology and say that Ω or $\mathcal{T}$ is the planar extension. Similarly, we say T has a *positive planar extension* when T has a planar extension with $0 < \mu(\Omega) < \infty$.

We will have occasional need to use a planar extension of T for which the invariant measure is Lebesgue measure. Let

$$(11) \quad \mathcal{Z}(x, y) = (x, y/(1 + xy))$$

and for each $M = \begin{pmatrix} a & b \\ c & d \end{pmatrix}$ as above, let $\widehat{\mathcal{T}}_M = \mathcal{Z} \circ \mathcal{T}_M \circ \mathcal{Z}^{-1}$. Then

$$(12) \quad \widehat{\mathcal{T}}_M(x, y) = (M \cdot x, (cx + d)^2 y - c(cx + d)).$$

Convention Throughout, we will allow ourselves the minor abuse of using adjectives such as injective, surjective and bijective to mean in each case *up to measure zero*, and thus similarly where we speak of disjointness and the like we again will assume the meaning being taken to include the proviso “up to measure zero” whenever reasonable.

2.4. Hyperbolic upper half-plane, geodesic flow, Fuchsian groups. The upper half-plane is $\mathbb{H} = \{z = x + iy \mid x, y \in \mathbb{R}, y > 0\}$, equipped with the hyperbolic metric on $\mathbb{H}$, whose element of arclength squares to be $ds^2 = (dx^2 + dy^2)/y^2$. The group $\mathrm{SL}(2, \mathbb{R})$ acts by fractional linear maps, which are isometries of $\mathbb{H}$. The geodesics of $\mathbb{H}$ are either vertical lines and semi-circles, whose naive extensions meet the boundary real line perpendicularly.

The action of $\mathrm{SL}_2(\mathbb{R})$ on $\mathbb{H}$ extends to an action on the unit tangent bundle $T^1\mathbb{H}$. This allows one to identify $T^1\mathbb{H}$ with $\mathrm{PSL}_2(\mathbb{R})$. The *geodesic flow* on the unit tangent bundle is an action of the real numbers: given a nonnegative real number t and a unit tangent vector u , the unit tangent vector $t \cdot u$ is obtained by following the unique geodesic passing through u in the positive direction for arclength t and taking the unit tangent vector to this oriented geodesic at the new basepoint. If $t < 0$, we follow the geodesic in the opposite direction. In terms of $\mathrm{PSL}_2(\mathbb{R})$, the geodesic flow for time t is given by sending $A \in \mathrm{PSL}_2(\mathbb{R})$ to Ag_t , where $g_t = \begin{pmatrix} e^{t/2} & 0 \\ 0 & e^{-t/2} \end{pmatrix}$.

A *Fuchsian group* Γ is a discrete (with respect to the natural topology) subgroup of $\mathrm{PSL}(2, \mathbb{R})$ and in particular, it is a subgroup acting properly discontinuously on $\mathbb{H}$. The quotient $\Gamma \backslash \mathbb{H}$ has a unit tangent bundle, given by equivalence classes (thus Γ -orbits) of unit tangent vectors of $\mathbb{H}$. The geodesic flow descends so that for $t \in \mathbb{R}$, $[A] \in \Gamma \backslash \mathrm{PSL}_2(\mathbb{R})$ is sent to $[Ag_t]$, where the square brackets here denote Γ -cosets represented by the given elements.

2.5. Volume of unit tangent bundles. The volume of the unit tangent bundle of a surface uniformized by a Fuchsian group Γ equals π times the hyperbolic area of a fundamental domain for Γ . Each of the $G_{m,n}$ with $2 \leq m$ and $3 \leq n$ is a hyperbolic triangle group. The fundamental domain is then the union of a hyperbolic triangle of angles $0, \pi/m, \pi/n$ with the reflection over one of its edges. The area of the fundamental domain is twice that of the triangle, but the area of a hyperbolic triangle equals π minus the sum of its angles. Therefore,

$$(13) \quad \text{vol}(T^1(G_{m,n} \setminus \mathbb{H})) = 2\pi (\pi - \pi/m - \pi/n) = 2\pi^2 \frac{mn - m - n}{mn}.$$

In particular, $\text{vol}(T^1(G_{2,n} \setminus \mathbb{H})) = \pi^2 (n - 2)/n$ and $\text{vol}(T^1(G_{3,n} \setminus \mathbb{H})) = 2\pi^2(2n - 3)/(3n)$.

2.6. Arnoux's method for cross sections to a geodesic flow.

2.6.1. *Cross sections to a measurable flow, Arnoux's transversal.* Let $(X, \mathcal{B}, \mu)$ be a measure space and Φ_t a measure preserving flow on X , that is $\Phi : X \times \mathbb{R} \rightarrow X$ is a measurable function such that for $\Phi_t(x) = \Phi(x, t)$, $\Phi_{s+t} = \Phi_s \circ \Phi_t$. Then $\Sigma \subset X$ is a *measurable cross section* for the flow Φ_t if: (1) the flow orbit of almost every point meets Σ ; (2) for almost every $x \in X$ the set of times t such that $\Phi_t(x) \in \Sigma$ is a discrete subset of $\mathbb{R}$; (3) for every Borel subset (in the subspace topology) $A \subset \Sigma$ and for every $\tau > 0$, *flow box* $A_{[0, \tau]} := \{\Phi_t(A) \mid A \in \Sigma, t \in [0, \tau]\}$ is μ -measurable.

The *return-time function* $r = r_\Sigma$ is $r(x) = \inf\{t > 0 : \Phi_t(x) \in \Sigma\}$ and the *return map* $R : \Sigma \rightarrow \Sigma$ is defined by $R(x) = \Phi_{r(x)}(x)$. The *induced measure* μ_Σ on Σ is defined from flow boxes: one sets $\mu_\Sigma(A) = \frac{1}{\tau} \mu(A_{[0, \tau]})$ for any $0 < \tau < \inf_{x \in A} \{r(x)\}$.

Convention: In all that follows, we write *cross section* to denote measurable cross-section.

A flow Φ_t is *ergodic* if for any invariant set either it or its complement is of measure zero. The fundamental result of Hopf states that if a Fuchsian group Γ is of finite covolume, then the geodesic flow is ergodic with respect to the natural measure on the unit tangent bundle of $\Gamma \setminus \mathbb{H}$; see his reprisal in [Ho]. A flow is *recurrent* if the Φ -orbit of almost every point meets any positive measure set infinitely often. By the Poincaré Recurrence Theorem, an ergodic flow on a finite measure space is recurrent. Given a cross-section, its first return-time transformation $(\Sigma, \mathcal{B}_\Sigma, \mu_\Sigma, R_\Sigma)$ is ergodic whenever the flow is.

There is a natural measure on the unit tangent bundle $T^1\mathbb{H}$: the *Liouville measure* is given as the product of the hyperbolic area measure on $\mathbb{H}$ with the length measure on the circle of unit vectors at any point. Liouville measure is (left- and right-) $\text{SL}_2(\mathbb{R})$ -invariant, and thus gives Haar measure on G . In particular, this measure is invariant for the geodesic flow. With standard normalizations, Liouville measure agrees with the Riemannian volume form. These both descend modulo any Fuchsian group Γ .

Arnoux, see say [AS], found an elementary manner to map Lebesgue planar extensions into the unit tangent bundle of appropriate surfaces so as to find cross sections to the geodesic flow. This often allows one to express an original interval map's system as a factor. For $(x, y) \in \mathbb{R}^2$, let

$$A(x, y) = \begin{pmatrix} x & xy - 1 \\ 1 & y \end{pmatrix}.$$

Arnoux's transversal, $\mathcal{A}$, is the projection to $\text{PSL}_2(\mathbb{R})$ of the set of all $A(x, y)$. The complement to $\{[Ag_t] \mid A \in \mathcal{A}, t \in \mathbb{R}\}$ is a null set for Liouville measure. Furthermore Liouville measure restricts to $\mathcal{A}$ to be a constant multiple of $dx dy dt$.

We repeat a definition and observation from [AS2].

Definition 2.1. For $M \in \text{SL}(2, \mathbb{R})$ and $x \in \mathbb{R}$ such that $M \cdot x \neq \infty$, let

$$\tau(M, x) := -2 \log |cx + d|$$

where (c, d) is the bottom row of M as usual.

Elementary calculation shows that τ induces on $\text{PSL}(2, \mathbb{R})$ a *cocycle*, in the following sense:

$$\tau(MN, x) = \tau(M, Nx) + \tau(N, x)$$

whenever all terms are defined and we choose each projective representative such that the corresponding $cx + d$ is positive.

For ease of legibility, set $t_0 = \tau(M, x)$. Then matrix multiplication shows

$$(14) \quad MAg_{t_0} = A(\widehat{\mathcal{T}}_M(x, y)).$$

Given a Fuchsian group Γ , we have the function

$$\begin{aligned} \mathcal{P} = \mathcal{P}_\Gamma : \mathbb{R}^2 &\rightarrow \Gamma \backslash \mathrm{PSL}(2, \mathbb{R}) \\ (x, y) &\mapsto [A(x, y)], \end{aligned}$$

where again square brackets denote cosets. Note that this is measure preserving when we use Lebesgue measure on $\mathbb{R}^2$ and the measure given by $dx dy$ inherited by the projection $\mathcal{A}_\Gamma$, of $\mathcal{A}$ to $\Gamma \backslash \mathrm{PSL}(2, \mathbb{R})$. When we use $\mathcal{P}$ we will often commit the abuse of writing A to represent the corresponding coset.

2.6.2. Arnoux's method. Now suppose that T is a piecewise (determinant one) Möbius interval map. We say that the *group associated to T* is the group Γ_T generated by the Möbius transformations of T ; thus, $\Gamma_T = \langle M_\beta, \beta \in \mathcal{B} \rangle \subset \mathrm{PSL}_2(\mathbb{R})$.

If T has a positive planar extension, we conjugate via $\mathcal{Z}$ of (11) to the Lebesgue planar extension of T . We then apply the measure preserving $\mathcal{P}$, with $\Gamma = \Gamma_T$, to $\mathcal{Z}(\Omega)$. When Γ_T is a Fuchsian group, this is a subset of $T^1(\Gamma_T \backslash \mathbb{H})$; using the fact that any Fuchsian group has countably many elements, in [AS2] it is shown that $\mathcal{P}$ is injective up to measure zero on $\mathcal{Z}(\Omega)$.

Arnoux's method is illustrated by the following result. When T is a piecewise determinant one Möbius interval map and $x \in \mathbb{I}$, we let $\tau(x) = \tau_T(x) = \tau(M, x)$ where $T(x) = M \cdot x$. The result here combines ([AS2] Theorem 5.4, Corollary 1, and Proposition 4).

Theorem 2.2. *[Arnoux's Method] Let T be a piecewise determinant one Möbius interval map with positive planar extension, and suppose that Γ_T is Fuchsian group of finite covolume. Then*

$$\Sigma = \mathcal{P}_\Gamma(\mathcal{Z}(\Omega))$$

is a cross section to the geodesic flow on $T^1(\Gamma_T \backslash \mathbb{H})$. Furthermore, the system defined by

$$\begin{aligned} \phi : \Sigma &\rightarrow \Sigma \\ [A(x, y)] &\mapsto [MA(x, y)g_{\tau(x)}], \end{aligned}$$

with M such that $T(x) = M \cdot x$, is an extension of $T : \mathbb{I} \rightarrow \mathbb{I}$. Moreover, ϕ agrees with the first return map of the geodesic flow to Σ if and only if T is ergodic, expansive, and with entropy satisfying $h(T)\mu(\Omega_f) = \mathrm{vol}(T^1(\Gamma_T \backslash \mathbb{H}))$. When this holds, the first return to Σ gives a natural extension to T .

The initial statement of the theorem is shown by considering (14) with the projection. The subset $\{[A(x, y)g_t] \mid A(x, y) \in \Sigma, 0 \leq t \leq \tau(x)\} \subset T^1(\Gamma_T \backslash \mathbb{H})$ is invariant under the geodesic flow; Hopf's result implies that this is all of $T^1(\Gamma_T \backslash \mathbb{H})$ up to measure zero. If ϕ agrees with the first return map, then the map to $\mathcal{A}$ is such that the flow is expansive in the x -direction and contracting in the y -direction — recall that the map of (12) preserves Lebesgue measure — and one can argue as in [AS2] that the first return system is indeed the natural extension. The ergodicity of the flow implies that ϕ is ergodic, it then follows that $\widehat{\mathcal{T}}, \mathcal{T}$ and hence T itself are ergodic. The veracity of the equation involving the entropy is also in [AS2]; in brief, $\tau(x)$ is simultaneously the arclength of the geodesic path following the flow line from $[A(x, y)]$ to its image under ϕ and the (piecewise form of the) integrand in Rohlin's formula (6). Recall that ν is the marginal measure of μ on Ω , and that both $\mathcal{Z}$ and $\mathcal{P}_\Gamma$ are measure preserving. That in this setting one has a natural extension can be shown, for instance, using the main result of [AS3].

2.7. Nakada's entropy formula for Rosen fractions. Recall that, for each $n \geq 3$, $\nu_n = 2 \cos \pi/n$; set $\mathbb{I}_n = [-\nu_n/2, \nu_n/2]$. Letting, as in [N2], $\lfloor x \rfloor_n = a\nu_n$ such that $x - a\nu_n \in \mathbb{I}_n$, Rosen's [R] continued fraction map on $\mathbb{I}_n$ is given by $f(x) = f_n(x) = |1/x| - \lfloor |1/x| \rfloor_n$ (with $x = 0$ fixed). Rosen introduced these maps to study the Hecke triangle Fuchsian groups, the $G_{2,n}$ in our notation.

Using a combination of hyperbolic geometrical considerations, Diophantine approximation, and ergodic theory, Nakada [N3] obtained the following. We slightly reformulate.

Theorem 2.3 (Nakada 2010). *Let f_n denote the Rosen continued fraction map of integral index $n \geq 3$ and Ω_n the planar domain as given in [BKS]. Then the product of the entropy of f_n times $\mu(\Omega_n)$ equals $1/2$ times the volume of the unit tangent bundle of the surface uniformized by the Hecke triangle Fuchsian group of index n , $G_{2,n}$.*

2.8. Setting of [CKS]. In [CKS] one has a detailed study of the interval maps $T_{3,n,\alpha}$. We tersely summarize their notation and results.

The parameter interval is naturally partitioned, $[0, 1] = \{0\} \cup (0, \gamma_n) \cup [\gamma_n, 1] \cup \{1\}$, where $\alpha < \gamma_n$ if $\forall x \in [\ell_0(\alpha), r_0(\alpha)]$ one has $T_{n,\alpha}(x) = A^k C \cdot x$ for some nonzero integer k . We call the set of $\alpha < \gamma_n$ the *small α* . For small α , the only possibly non-full cylinders are the leftmost cylinder (whose left endpoint is $\ell_0(\alpha)$) and the rightmost (with right endpoint $r_0(\alpha)$).

For $x \in [\ell_0(\alpha), r_0(\alpha)]$ we call $d_\alpha(x) = k$ the (simplified) α -digit of x , and say that x lies in the cylinder $\Delta_\alpha(k)$. To each x we then associate the sequence of its α -digits. The T_α -orbits of the endpoints $\ell_0(\alpha), r_0(\alpha)$ are of extreme importance in our discussions, one writes $\underline{b}_{[1,\infty)}^\alpha$ for the α -digits of $\ell_0(\alpha)$ and replaces the lower bar by an upper bar, $\bar{b}_{[1,\infty)}^\alpha$ for $r_0(\alpha)$. We also label entries in these orbits by $\ell_0, \ell_1, \dots$ and $r_0, r_1, \dots$.

2.8.1. Expansions given words, synchronization intervals, tree of words. For small α , we define parameter intervals on which an initial portion of the expansion of the $r_0(\alpha)$ are fixed. That is, there is a common prefix of the $\bar{d}_{[1,\infty)}^\alpha$. Given a word $v = c_1 d_1 \cdots d_{s-1} c_s$ of positive integers, let $\bar{S}(v) = \sum_{i=1}^s c_i + \sum_{j=1}^{s-1} d_j$ and for each $k \in \mathbb{N}$ let

$$(15) \quad \bar{d}(k, v) = k^{c_1}, (k+1)^{d_1}, \dots, (k+1)^{d_{s-1}}, k^{c_s}, \quad \text{when } v = c_1 d_1 \cdots d_{s-1} c_s.$$

The corresponding subinterval of parameters is

$$(16) \quad I_{k,v} = \{\alpha \mid \bar{d}_{[1,\bar{S}(v)]}^\alpha = \bar{d}(k, v)\};$$

in words, this is the set of α for which the α -expansion (in simplified digits) of $r_0(\alpha)$ has $\bar{d}(k, v)$ as a prefix. Thus, setting

$$(17) \quad R_{k,v} = (A^k C)^{c_s} (A^{k+1} C)^{d_{s-1}} (A^k C)^{c_{s-1}} \cdots (A^{k+1} C)^{d_1} (A^k C)^{c_1},$$

for $\alpha \in I_{k,v}$ one has $T_\alpha^{\bar{S}(v)}(r_0(\alpha)) = R_{k,v} \cdot r_0(\alpha)$. The left endpoint of $I_{k,v}$ is denoted $\zeta_{k,v}$, one finds that $R_{k,v} \cdot r_0(\zeta_{k,v}) = \ell_0(\zeta_{k,v})$. We define

$$(18) \quad J_{k,v} = [\zeta_{k,v}, \eta_{k,v}),$$

where with $L_{k,v} = C^{-1} A C R_{k,v}$ we have $L_{k,v} \cdot r_0(\eta_{k,v}) = r_0(\eta_{k,v})$. Confer ([CKS], Figure 4.2).

A main result of [CKS] is that each $J_{k,v}$ is what is now usually called a *matching interval*: for all α in the interior of $J_{k,v}$, the T_α -orbits of $\ell_0(\alpha)$ and of $r_0(\alpha)$ meet and do so in a common fashion. We furthermore showed that the complement in $(0, \gamma_n)$ of the union of the $J_{k,v}$ is of measure zero. Key to this was determining a maximal common prefix of the $\bar{d}_{[1,\infty)}^\alpha$ for $\alpha \in J_{k,v}$. For that, we used $w = w_n = (-1)^{n-2}, -2, (-1)^{n-3}, -2$ and for k fixed, we let $\mathcal{C} = \mathcal{C}_k =$

$(-1)^{n-3}, -2, w^{k-1}$ and $\mathcal{D} = C_{k+1}$, and defined $\underline{d}(k, v) = w^k, \mathcal{C}^{c_1-1} \mathcal{D}^{d_1} \dots \mathcal{D}^{d_{s-1}} \mathcal{C}^{c_s}, (-1)^{n-2}$. This is the common prefix, and also

$$(19) \quad \underline{d}_{[1, \infty)}^{n_{k,v}} = \overline{w^k, \mathcal{C}^{c_1-1} \mathcal{D}^{d_1} \dots \mathcal{D}^{d_{s-1}} \mathcal{C}^{c_s}, (-1)^{n-3}, -2},$$

where the overline indicates a period. We denote the length as a word in $\{-1, -2\}$ of $\underline{d}(k, v)$ by $\underline{S}(k, v)$. There is an expression for $L_{k,v}$ related to $\underline{d}(k, v)$ in a manner similar to how (17) relates $R_{k,v}$ to $\bar{d}(k, v)$. The aforementioned synchronization is of the form

$$(20) \quad T_\alpha^{\bar{S}(v)+1}(r_0(\alpha)) = T_\alpha^{\underline{S}(v)+1}(\ell_0(\alpha)).$$

The words v that we use form a tree, $\mathcal{V}$. For this, we first define v' such that any $\bar{d}_{[1, \infty)}^{n_{k,v}}$ equals $\bar{d}(k, v(v')^\infty)$, where the latter is a slight abuse of notation. The following is [[CKS], Definition 4.8].

Definition 2.4. For each $s > 1$ and each word $v = c_1 d_1 \dots c_{s-1} d_{s-1} c_s$, define

$$v' = \begin{cases} 1(c_1 - 1)d_1 c_2 \dots c_{s-1} d_{s-1} c_s & \text{if } c_1 \neq 1, \\ (d_1 + 1)c_2 \dots c_{s-1} d_{s-1} c_s & \text{otherwise.} \end{cases}$$

We interpret this also to mean that when $v = c$ with $c > 1$ then $v' = 1(c - 1)$, and when $v = 1$ then $v' = 1$.

The following is [[CKS], Definition 4.10].

Definition 2.5. Set $\Theta_{-1}(c_1) = c_1 + 1$ and $\Theta_q(1) = 1q1$ for $q \geq 1$. For $c > 1$, set $\Theta_q(c) = c[1(c - 1)]^q 1c$ for any $q \geq 0$. (To avoid double labeling and also to stay within our desired language, $\Theta_0(1)$ is undefined; note that $\Theta_1(1) = 111$, compare with $\Theta_0(c) = c1c$ for $c > 1$.)

We now recursively define values of the operators Θ_q . Suppose $v = \Theta_p(u) = uv''$ for some $p \geq 0$ and some suffix v'' . Then define for any $q \geq 0$

$$\Theta_q(v) = v(v')^q v''.$$

Definition 2.6. Let $\mathcal{V}$ denote the set of all words obtainable from $v = 1$ by finite sequences of applications of the various Θ_q .

There is a type of self-similarity of $\mathcal{V}$ which allows the explicit definition of the *derived words* operator $\mathcal{D}$ such that (again for general v) $\mathcal{D} \circ \Theta_q(v) = \Theta_q \circ \mathcal{D}(v)$. In general $\mathcal{D}$ decreases the length of words while preserving various properties, and hence assists in induction proofs. For example, we applied it to prove that every $v \in \mathcal{V}$ is a palindrome. For more on the tree, see ([CKS], §4.2).

2.9. First expansive power maps. We recall the following directly from [CKS3].

Definition 2.7. Let T be a function on an interval I , with an invariant probability measure ν which is equivalent to Lebesgue measure. Suppose that T is eventually expansive, thus there is some $r \in \mathbb{N}$ such that T^r is expansive. By the Well Ordering Principle of the integers, for each $x \in I$ there is a least $\ell(x) \in \mathbb{N}$ such that $T^{\ell(x)}(x)$ is expansive in the sense that the derivative here is greater than 1 in absolute value. We define the *first pointwise expansive power* of T as the map $U : I \rightarrow I$ sending each x to $T^{\ell(x)}(x)$.

Of course, $\ell(x) \leq r$ for all $x \in I$ and also $U(x) = T(x)$ for all $x \in E$, where E is the maximal subset of I on which T itself is expansive. For each $k \leq r$, let $E_k \subset I$ be the set on which $\ell(x) = k$. Of course, the E_k give a finite partition of I . We have that $|(T^k)'(x)| > 1$ for all $x \in E_k$; by the Chain Rule, $|(T^k)'(x)| = \prod_{i=0}^{k-1} |T'(T^i(x))|$. The *minimality* of $\ell(x)$ hence implies both $T(E_k) \subseteq \cup_{i=1}^{k-1} E_i$ for $k \geq 2$ and $T^{k-1}(E_k) \subset E_1$ for all k .

Note that in the setting where I is an interval and T is given piecewise by Möbius transformations, then also U is so given.

3. RESULTS FOR THE SETTING OF $m = 2$

Besides the aforementioned Theorem 2, we show the following results specific to the $m = 2$ setting.

Theorem 3. *For each $n \geq 3$, the function assigning to $\alpha \in (0, 1)$ the measure theoretic entropy of $T_{2,n,\alpha}$ is continuous and is invariant under $\alpha \mapsto 1 - \alpha$.*

For each $n \geq 3$ and $\alpha \in [0, 1]$ let $\Omega_{2,n,\alpha}$ be as defined in Subsections § 3.7, 3.9 and 3.10, according as α : lies in the interior of a matching interval; is any other value in $(0, 1)$; is in $\{0, 1\}$.

Theorem 4. *For each $n \geq 3$ and $\alpha \in [0, 1]$, the map $(x, y) \mapsto (-x, -y)$ sends $\Omega_{2,n,\alpha}$ to $\Omega_{2,n,1-\alpha}$.*

Theorem 5. *For all $n \geq 3$ and $\alpha \in [0, 1]$ we have*

$$\int_{\Omega_{2,n,\alpha}} \log |T'_{2,n,\alpha}(x)| d\mu = \text{vol}(T^1(G_{2,n} \setminus \mathbb{H})).$$

For $\alpha \in (0, 1)$, the measure theoretic entropy of $T_{2,n,\alpha}$ satisfies

$$h(T_{2,n,\alpha})\mu(\Omega_{2,n,\alpha}) = \text{vol}(T^1(G_{2,n} \setminus \mathbb{H})).$$

In this section, we first consider in § 3.1 the basic symmetry $\alpha \mapsto 1 - \alpha$ between interval maps and also between their associated planar maps. This symmetry naturally points to $\alpha = 1/2$ as a special case; in § 3.2 we prove the second statement of Theorem 5 in the case of $\alpha = 1/2$, a result used without proof in [AS2] and which is indeed virtually a direct result of Nakada's Theorem 2.3. For each $n \geq 3$, the exact value of the μ -mass of $\Omega_{2,n,1/2}$ is twice that of the planar extension region given by [BKS]. Due to the result of § 3.2, Birkhoff approximations for the entropy of certain $T_{2,n,1/2}$ made us suspicious of their μ -mass value in the odd n case. In § 3.3 we correct their formula.

The remainder of the section, other than § 3.8, shows that the techniques and results of [CKS, CKS3, CKS2] hold in the $n = 2$ setting. In § 3.8 we give the main results underpinning Theorem 2. That is, we determine the behavior of $\alpha \mapsto \Omega_{2,n,\alpha}$ along matching intervals, determine the middle matching interval along which the function is constant, as well as the maximal value.

In § 3.9 we complete the proof of all but the first statement of Theorem 2 by showing the continuity of $\alpha \mapsto \mu(\Omega_{2,n,\alpha})$ on all of $(0, 1)$. In § 3.10, we determine $\Omega_{2,n,0}$ and $\Omega_{2,n,1}$ for all $n \geq 3$. Combined with the earlier constructions of the various $\Omega_{2,n,\alpha}$, this completes the proof of Theorem 4. In an initial part of § 3.11, we use the explicit two dimensional regions $\Omega_{2,n,\alpha}$ and 'word processing' to allow a use of Abramov's Formula, showing continuity of $\alpha \mapsto h(T_{2,n,\alpha})$ on all of $(0, 1)$; this completes both the proof of Theorem 2, and due to the symmetry, also that of Theorem 3. In the later part of § 3.11, we consider Rohlin integrals to complete the proof of Theorem 5.

3.1. Symmetry under $\alpha \mapsto 1 - \alpha$. The basic observation that for $x \notin \mathbb{Z}$ one has $\lfloor -x \rfloor = -1 - \lfloor -x \rfloor$ leads to the following.

3.1.1. *Interval maps.* If $T_{2,n,\alpha}(x) \neq \ell_0(\alpha)$, one finds $T_{2,n,1-\alpha}(-x) = -T_{2,n,\alpha}(x)$. Furthermore, if $T_{2,n,\alpha}(x) = \ell_0(\alpha) = (\alpha - 1)t$ then $T_{2,n,1-\alpha}(-x) = \ell_0(1 - \alpha) = -\alpha t$.

For more detail, when $T_{2,n,\alpha}(x) \neq \ell_0(\alpha)$, one has $\frac{1}{tx} + 1 - \alpha \notin \mathbb{Z}$ and hence

$$\begin{aligned} -T_{2,n,\alpha}(x) &= \frac{1}{x} + \lfloor \frac{-1}{tx} + 1 - \alpha \rfloor t && =: \frac{1}{x} + kt \\ &= \frac{1}{x} + (1 + \lfloor \frac{-1}{tx} - \alpha \rfloor) t \\ &= \frac{1}{x} + (1 - 1 - \lfloor \frac{1}{tx} + \alpha \rfloor) t \\ &= \frac{1}{x} - \lfloor \frac{1}{tx} + \alpha \rfloor t \\ &= \frac{1}{x} - \lfloor \frac{1}{tx} + 1 - (1 - \alpha) \rfloor t && =: \frac{1}{x} - k't \\ &= T_{2,n,1-\alpha}(-x). \end{aligned}$$

Thus, from $k' = -k$ in the above, one finds that for any x in the interior of $\mathbb{I}_{2,n,\alpha}$ with $T_{2,n,\alpha}(x) = A^k C \cdot x \neq \ell_0(\alpha)$, one has $-x \in \mathbb{I}_{2,n,1-\alpha}$ and $T_{2,n,1-\alpha}(-x) = -T_{2,n,\alpha}(x) = C(A^k C)^{-1} C \cdot (-x)$.

3.1.2. *Planar maps.* From the previous, for x as above and any y -value, since when $m = 2$ we have $R = C$,

$$\begin{aligned} \mathcal{T}_{2,n,1-\alpha}(-x, -y) &= (-A^k C \cdot x, RC(A^k C)^{-1} CR^{-1} \cdot -y) \\ &= (-A^k C \cdot x, (A^k C)^{-1} \cdot -y) \\ &= (-A^k C \cdot x, \frac{1}{kt + y}) \\ &= -(A^k C \cdot x, RA^k C R^{-1} \cdot y) \\ &= -\mathcal{T}_{2,n,\alpha}(x, y). \end{aligned}$$

Thus, since $d\mu = (1 + xy)^{-2} dx dy$ is preserved by

$$(21) \quad \mathcal{S} : (x, y) \mapsto (-x, -y),$$

the bijectivity up to μ -null sets of $\mathcal{T}_{2,n,\alpha}$ on a domain, such as is shown for $\Omega_{2,n,\alpha}$ below, is equivalent to that of $\mathcal{T}_{2,n,1-\alpha}$ on $\mathcal{S}(\Omega_{2,n,\alpha})$. Thus, the symmetry announced in Theorem 4 will hold.

3.2. Entropy of symmetric Rosen fractions. Given the above symmetry, the value $\alpha = 1/2$ has heightened interest. We fill in missing details from [AS2]. There, a family of interval maps called the “symmetric Rosen maps” are introduced and studied. In fact, in our notation the symmetric Rosen map of index n is $g = g_n = T_{2,n,1/2}$. For each $n \geq 3$, their definition begins with setting $\lambda_n = 2 \cos \pi/n$ and $\mathbb{I}_n = [-\lambda_n/2, \lambda_n/2]$. Letting, as in [N2], $\lfloor x \rfloor_n = a\lambda_n$ such that $x - a\lambda_n \in \mathbb{I}_n$, Rosen’s [R] continued fraction map on $\mathbb{I}_n$ is given by $f(x) = f_n(x) = |1/x| - \lfloor |1/x| \rfloor_n$ (with $x = 0$ fixed). Similarly, [AS2] defines the symmetric Rosen map $g(x) = g_n(x) = -1/x - \lfloor -1/x \rfloor_n$. In [AS2] this function is denoted by $h(x)$, and the indexing variable is q instead of n . Since $\lambda_n = t_{2,n}$, one easily verifies that $g_n = T_{2,n,1/2}$.

Of course, $f(x) = g(x)$ if $x \leq 0$. On the other hand, for $x > 0$, we have $f(x) = g(-x)$. As noted in [AS2] $g(-x) = -g(x)$ unless $g(x) = -\lambda_n/2$. Thus, for almost all $x > 0$, $g'(x) = (-g(-x))' = g'(-x) = f'(-x)$. Hence, $|g'(x)| = |f'(x)|$ on $\mathbb{I}_n$, up to a null set.

A planar natural extension Ω_n is given for each f_n in [BKS]; for each, $\forall (x, y) \in \Omega_n$, $y \geq 0$. Let $\mathcal{G} = \mathcal{G}_n$ be the planar map associated to g in our usual manner. We now justify the implicit claim of [AS2] that $\mathcal{G}$ is bijective up to μ -null sets on the union $\Omega_n \cup \mathcal{S}(\Omega_n)$. Let $\mathcal{F}$ denote the

two dimensional map related to $f = f_n$ in our usual manner. We have $\mathcal{G}(x, y) = \mathcal{F}(x, y)$ for $x \leq 0$ and $(x, y) \in \Omega_n$ and $\mathcal{F}(x, y) = \mathcal{G}(-x, -y)$ for $x > 0$ and $(x, y) \in \Omega_n$. Therefore, up to null sets, $\mathcal{G}$ sends $\{x \leq 0\} \cap \Omega_n$ bijectively to Ω_n . The symmetry $\mathcal{G} \circ \mathcal{S}(x, y) = -\mathcal{G}(x, y)$ now shows that $\mathcal{G}$ is bijective up to null sets on $\Omega_n \cup \mathcal{S}(\Omega_n)$. Note that in our notation, $\mathcal{G} = T_{2,n,1/2}$ and one can verify that $\Omega_n \cup \mathcal{S}(\Omega_n) = \Omega_{2,n,1/2}$.

Since the measure μ is invariant under the symmetry $\mathcal{S}$, we have $\mu(\Omega_n \cup \mathcal{S}(\Omega_n)) = 2\mu(\Omega_n)$. Applying the Rohlin formula (6), the entropy of the symmetric Rosen map $g_n(x)$ satisfies

$$\begin{aligned} h(g_n) \mu(\Omega_n \cup \mathcal{S}(\Omega_n)) &= \int_{\Omega_n \cup \mathcal{S}(\Omega_n)} \log |g'(x)| d\mu \\ &= 2 \int_{\Omega_n} \log |g'(x)| d\mu \\ &= 2 \int_{\Omega_n} \log |f'(x)| d\mu. \end{aligned}$$

Thus, Nakada's result, Theorem 2.3, then gives

$$(22) \quad h(T_{2,n,1/2}) \mu(\Omega_{2,n,1/2}) = \text{vol } T^1(G_{2,n} \backslash \mathbb{H}).$$

3.3. Measure of Ω_n for Rosen fractions: Correcting an (typographical?) error in [BKS]. There is little doubt but that the author of this paper is responsible for an error in [BKS]. In short, [[BKS], Lemma 3.4] is wrong as stated. The error looks to be caused by what seems to be a typographical error in the formula for the μ -mass of the planar extension there at hand. We correct this latter in the following.

The μ -mass of the planar extension for the Rosen interval map of odd index $n \geq 3$ is given by

$$(23) \quad \mu(\Omega_n) = \ln \left(\frac{1 + R_n}{2 \sin \frac{\pi}{2n}} \right),$$

where R_n is the positive solution of $x^2 + (2 - 2 \cos \pi/n)x - 1 = 0$. Note that this expression for $\mu(\Omega_n)$ also holds for $n = 3$, a case not considered in [BKS], where the Rosen continued fraction is the standard 'nearest integer' continued fraction.

The denominator in the corrected expression is the multiplicative inverse of the [BKS] quantity $B_{h+1} := \sin \frac{(h+1)\pi}{n} / \sin \frac{\pi}{n}$ with $n = 2h + 3$. The flaw in [BKS] is either simply a dropping of this factor, or is based upon an error in the proof of [[BKS], Lemma 3.4] where a certain denominator is written as being B_{j+2} but should be B_{j+1} . With that corrected, a telescoping product has cancellation that leaves B_{h+1} in place of simply 1.

The formula can be compared with the (correct!) formula when n is even for $\mu(\Omega_n)$ given by [BKS]: $\mu(\Omega_{\text{even } n}) = \ln \left[\frac{1 + \cos \pi/n}{\sin \pi/n} \right]$, see Figure 2.

3.4. Matching relation. Given k and $v \in \mathcal{V}$, we have $R_{k,v}$ as in (17). We now follow [CKS] to define $L_{k,v}$. With $m = 2$, the word w above becomes $w = (-1)^{n-3}, -2$. Correspondingly we let

$$(24) \quad W = W_{2,n} = A^{-2}C(A^{-1}C)^{n-3}$$

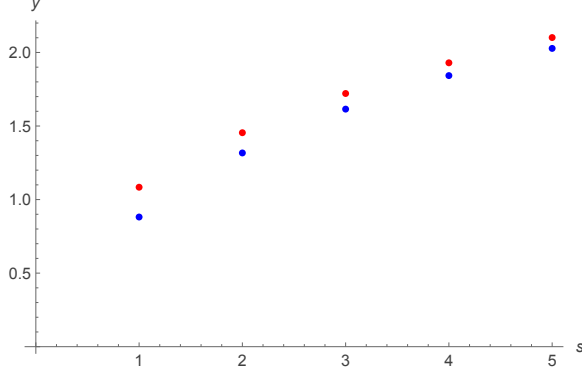


FIGURE 2. The μ -mass of the Rosen planar extension Ω_{2s+3} of [BKS] (in red) is comparable to that of Ω_{2s+2} (in blue), with both going to infinity as $s \rightarrow \infty$. See § 3.3.

and consider the following, each a word equal to $W^{k-1}(A^{-1}C)^{-1}$,

$$\tilde{C}_k = \begin{cases} W^{k-2}A^{-2}C(A^{-1}C)^{n-4} & \text{if } n > 3, k > 1 \\ W^{k-2}A^{-1} & \text{if } n = 3, k > 1 \\ (A^{-1}C)^{-1} & \text{if } n > 3, k = 1. \end{cases}$$

For ease, let

$$\tilde{C} = \tilde{C}_k \text{ and } \tilde{D} = \tilde{C}_{k+1}$$

and define

$$(25) \quad L_{k,v} = (A^{-1}C)^{n-3} \tilde{C}^{c_s} \tilde{D}^{d_{s-1}} \dots \tilde{D}^{d_1} \tilde{C}^{c_1-1} W^{k-1} A^{-1},$$

where when $k = 1$ we simplify adjacent occurrences of powers of $A^{-1}C$, which as we will see, in every case leads to only nonnegative powers appearing.

Example 3.1. Consider $n = 3$ with $(k, v) = (2, p)$. Thus, $R_{k,v} = (A^2C)^p$ and $L_{2,p} = (A^{-1}C)^{n-3} \tilde{C}^{p-1} W^{k-1} A^{-1} = (A^{-1})^{p-1} W A^{-1} = (A^{-1})^p A^{-2} C A^{-1} = A^{-(p+2)} C A^{-1}$. Similarly, still with $n = 3$, for $(k, v) = (2, 212)$ we find

$$\begin{aligned} L_{2,212} &= \tilde{C}^2 \tilde{D} \tilde{C} W A^{-1} \\ &= (A^{-1})^2 A^{-2} C A^{-1} A^{-1} A^{-2} C A^{-1} \\ &= A^{-4} C A^{-4} C A^{-1} \end{aligned}$$

Let's now try $v = 121$.

$$\begin{aligned} L_{2,121} &= \tilde{C} \tilde{D}^2 W A^{-1} \\ &= A^{-1} (A^{-2} C A^{-1}) (A^{-2} C A^{-1}) A^{-2} C A^{-1} \\ &= (A^{-3} C)^3 A^{-1}. \end{aligned}$$

Let's now try the longer $v = 12121$.

$$\begin{aligned} L_{2,12121} &= \tilde{C} \tilde{D}^2 \tilde{C} \tilde{D}^2 W A^{-1} \\ &= A^{-1} (A^{-2} C A^{-1}) (A^{-2} C A^{-1}) A^{-1} (A^{-2} C A^{-1}) (A^{-2} C A^{-1}) A^{-2} C A^{-1} \\ &= (A^{-3} C)^2 A^{-4} C (A^{-3} C)^2 A^{-1}. \end{aligned}$$

Example 3.2. Consider $n = 7$ with $(k, v) = (1, 2)$. We have $R_{k,v} = (AC)^2$ and $L_{1,2} = (A^{-1}C)^{n-3} \tilde{\mathcal{C}}^{c_1-1} W^{k-1} A^{-1} = (A^{-1}C)^4 (A^{-1}C)^{-1} A^{-1} = (A^{-1}C)^3 A^{-1}$. Similarly, still with $n = 7$, for $(k, v) = (1, 212)$ we find

$$\begin{aligned} L_{1,212} &= (A^{-1}C)^{n-3} \tilde{\mathcal{C}}^2 \tilde{\mathcal{D}} \tilde{\mathcal{C}} A^{-1} \\ &= (A^{-1}C)^4 (A^{-1}C)^{-2} A^{-2} C (A^{-1}C)^{n-4} (A^{-1}C)^{-1} A^{-1} \\ &= (A^{-1}C)^2 A^{-2} C (A^{-1}C)^2 A^{-1} \end{aligned}$$

Thus, with $n = 7$ we have $\underline{S}(1, 212) = \overline{S}(212)$.

Compare the following with [[CKS], Lemma 5.1]. Fix n and take $R_{k,v}$ as in (17).

Lemma 3.3. *The following identities in $G_{2,n}$ hold.*

- (i) $A^k C A = C A W^k$ for each $k \in \mathbb{Z}$;
- (ii) $(C^{-1} A C)(A^k C)^a (C^{-1} A C)^{-1} = (A^{-1} C)^{n-3} \tilde{\mathcal{C}}_k^a (A^{-1} C)^{-(n-3)}$ for all $a, k \in \mathbb{Z}$;
- (iii) $\tilde{\mathcal{C}}_k (A^{-1} C)^{-(n-3)} (C^{-1} A C) = W^{k-1} A^{-1}$;
- (iv) $(C^{-1} A C) R_{k,v} = L_{k,v}$.

Proof. We have $C A W (C A)^{-1} = C A A^{-2} C (A^{-1} C)^{n-3} (C A)^{-1} = (C A^{-1})^{n-1} C^{-1} = A C^{-1} C^{-1} = A$, as C has projective order two. Thus for any k , we have $C A W^k (C A)^{-1} = A^k$ from which (i) holds.

Now, $(C^{-1} A C) A^k C = (C^{-1} A C) C A W^k A^{-1} = C^{-1} A^2 W^k A^{-1} = (A^{-1} C)^{n-3} W^{k-1} A^{-1}$. Since $\tilde{\mathcal{C}}_k$ equals $W^{k-1} (A^{-1} C)^{-1}$, one has $(C^{-1} A C) A^k C = (A^{-1} C)^{n-3} \tilde{\mathcal{C}}_k (A^{-1} C) A^{-1}$. Since $(A^{-1} C)^{n-3} = (A^{-1} C)^3 = (A^{-1} C) A^{-1} (C^{-1} A C)^{-1}$, the second statement holds.

Since $(A^{-1} C)^{-(n-3)} (C^{-1} A C) = (A^{-1} C)^3 C^{-1} A C = (A^{-1} C)^2 C = A^{-1} C A^{-1}$, we find $W^{k-1} A^{-1} = W^{k-1} (A^{-1} C)^{-1} A^{-1} C A^{-1} = \tilde{\mathcal{C}}_k (A^{-1} C)^{-(n-3)} (C^{-1} A C)$. Thus, the third statement holds.

From (ii) and (iii),

$$(C^{-1} A C)(A^k C)^{c_1} = (A^{-1} C)^{n-3} \tilde{\mathcal{C}}_k^{c_1-1} \tilde{\mathcal{C}}_k (A^{-1} C)^{-(n-3)} (C^{-1} A C) = (A^{-1} C)^{n-3} \tilde{\mathcal{C}}_k^{c_1-1} W^{k-1} A^{-1}.$$

Thus (iv) holds when $s = 1$. When $s > 1$, repeated applications of (ii) give

$$(C^{-1} A C) R_{k,v} = (A^{-1} C)^{n-3} \tilde{\mathcal{C}}^{c_s} \tilde{\mathcal{D}}^{d_{s-1}} \dots \tilde{\mathcal{D}}^{d_1} (A^{-1} C)^{-(n-3)} (C^{-1} A C) \tilde{\mathcal{C}}^{c_1}.$$

Thus, the previous displayed equation applies to show that here also (iv) holds. $\square$

3.5. Initial digits of right endpoints. We recycle the notation of § 2.8 to the setting here of $m = 2$. Fix $n \geq 3$ and set $t = t_{2,n}$. Since $AC \cdot x = t - 1/x$ (here A now denotes $A_{2,n}$ and C denotes C_2), an elementary calculation shows that given $\zeta_{1,1} = 1 - \sqrt{1 - 1/t^2}$. The largest value of x we contemplate is $x = t$; since $AC \cdot t = t^2 - 1/t$ is less than or equal to t for all $n \geq 3$, we find that $I_{1,1} = [1 - \sqrt{1 - 1/t^2}, 1]$ for $n > 3$. Note that this value, $\zeta_{1,1}$ is less than $1/2$. When $n = 3$ one finds that exactly when $\alpha = 1$ does $T_{2,3,\alpha}(x)$ ever equal $AC \cdot x$ and in fact this equality holds only for $x = 1$. A similar direct calculation shows that for $n = 3$, we have $I_{2,1}$ extending from $\zeta_{2,1}$ to 1 and that $\zeta_{2,1}$ is less than $1/2$. In all cases, moving to the left from this rightmost subinterval of the parameters, one finds for all possible greater values of k that $r_1(\alpha) = A^k C \cdot r_0(\alpha)$ does appear. Thus, the parameter intervals $0 < \alpha < 1$ are partitioned by the subintervals $I_{k,1}$ indexed by the positive k , although (by a minor abuse) $I_{1,1} = \emptyset$ when $m = 3$.

We now say more about the orbit of points in $I_{1,1}$ when $n > 3$. Upon setting $m = 2$ the element of [[CKS], Eq. 3.1] becomes $U = A^2 C (AC)^{n-3}$; the final lines of the proof of

[[CKS], Prop. 3.1] show that $U \cdot t = t$ (with an argument that also holds in our case of $m = 2$). From this, one can show that for $n > 3$ when $k = 1$ the only possible v are limited to have $c_i \leq n - 3$. Indeed, for any $x < 0$ we have $AC \cdot x = t - 1/x > t \geq r_0(\alpha)$ and hence already AC never appears for any negative x . A direct calculation with $0 < x < t$ gives $(AC)^{n-2} \cdot x = (CA^{-1})^2 \cdot x = -1/(-t-1/(-t+x)) = (x-t)/((x-t)t+1)$ and this is easily shown to be less than $x - t$. This last signals that $(AC)^{n-2}$ is inadmissible. That is, when $n > 3$ the digit $k = 1$ can occur, but the only admissible $R_{1,v}, v = c_1 d_1 \dots d_{s-1} c_s$ are such that $c_i < n - 2$ for all i .

3.6. Matching intervals. The description of the matching intervals when $m = 3$ of [[CKS], Theorem 4.12] needs only mild adjustment to give the following. We use the notation of § 2.8.1, with the interpretation throughout that elements are determined in terms of $(m, n) = (2, n)$.

Theorem 6. *Fix $m = 2$ and $n \geq 3$. Then*

$$(0, 1) = \bigcup_{k=1}^{\infty} I_{k,1}.$$

If $n = 3$ then $I_{1,1} = \emptyset$. Furthermore, for each $k \in \mathbb{N}$ and each $v \in \mathcal{V}$, one has

$$I_{k,v} = J_{k,v} \cup \bigcup_{q=q_0}^{\infty} I_{k,\Theta_q(v)},$$

where $q_0 = 0$ unless $v = c_1$, in which case $q_0 = -1$. Moreover, when $n > 3$ one has $I_{1,v} = \emptyset$ for all $v = c_1 \dots$ with $c_1 > n - 3$.

The proof of the analog of the theorem in the case of $m = 3$ is given in [[CKS], Subsections 4.3 and 4.4]. This proof succeeds *mutatis mutandis* for the main cases above. The special restrictions in our setting of $m = 2$ are a consequence of the discussion in the previous subsection.

3.7. Planar extensions along matching intervals.

3.7.1. Planar extensions $\Omega_{2,n,\alpha}$ for α interior to a matching interval. In [[CKS3], Table 1] one is pointed to the appropriate definitions of what we denote as $\Omega_{3,n,\alpha}$. Those definitions for their case of ‘small α ’ apply directly to the case of $m = 2$, giving the various $\Omega_{2,n,\alpha}$ with $\alpha \in (0, 1)$. A main reason that the constructions and arguments of [CKS3] succeed even when $m = 2$ is that the key lemma, [[CKS3], Lemma 14] is based on the fact that matrices $A^p C, p \in \mathbb{Z}$ are of the shape $\begin{pmatrix} a & b \\ -b & 0 \end{pmatrix}$. Since that is also true when $m = 2$, for real x , we have

$$(26) \quad (RA^p C R^{-1})^{-1} \cdot (-x) = -(A^p C \cdot x).$$

With this, we can pass from understanding of sequences of digits to actions of powers of our $\mathcal{T}_{2,n,\alpha}$.

Fixing (k, v) , recall that $\underline{S} = \underline{S}(k, v)$ denotes the length of $\underline{d}(k, v)$ as a word in $\{-1, -2\}$. For α interior to a matching interval $J_{k,v}$, let

$$\Omega^+ = ([\underline{\ell}_{\underline{S}+1}(\alpha), r_0(\alpha)) \times [0, -\underline{\ell}_{\underline{S}-i_{\underline{S}+1}}(\zeta_{k,v})] \cup \bigcup_{a=1}^{\underline{S}+1} ([\ell_{i_a}(\alpha), \ell_{i_{a+1}}(\alpha)) \times [0, -\underline{\ell}_{\underline{S}-i_a}(\zeta_{k,v})]),$$

where $\ell_{i_1}(\alpha), \dots, \ell_{i_{\underline{S}+1}}(\alpha)$ gives the initial $T_{2,n,\alpha}$ -orbit of $\ell_0(\alpha)$ in increasing order as real numbers, and similarly for the $\ell_{i_a}(\zeta_{k,v})$. It is helpful to define the values $\tau(i)$ by $\tau(j) = a$ exactly when $i_a = j$. Thus, $y_{\tau(i)} = -\underline{\ell}_{\underline{S}-i}(\zeta_{k,v})$. That the y_i are increasing is shown in [[CKS3], Lemma 35]. Hence, $y_{\underline{S}+1} = -\ell_0(\zeta_{k,v})$.

Similarly, let

$$\Omega^- = ([\ell_0(\alpha), r_{\bar{S}}(\alpha)] \times [-r_{j_{-1}}(\eta_{k,v}), 0]) \cup \bigcup_{b=-1}^{-\bar{S}} ([r_{j_{b-1}}(\alpha), r_{j_b}(\alpha)] \times [-r_{\bar{S}+b}(\eta_{k,v}), 0]),$$

where now $r_{j_{-1}}(\alpha), r_{j_{-2}}(\alpha), \dots, r_{j_{-\bar{S}-1}}(\alpha)$ denotes the initial $T_{2,n,\alpha}$ -orbit of $r_0(\alpha)$ in decreasing order, and similarly for the $r_{j_b}(\eta_{k,v})$ and $\bar{S} = \bar{S}(v)$ is the length of $\bar{d}(k, v)$ as a word in $\{k, k+1\}$. Again a result in [CKS3] shows that $y_{-1} > y_{-2} > \dots > y_{-\bar{S}-1}$. Note that $y_{-\bar{S}-1} = -r_0(\eta_{k,v})$.

The arguments of [[CKS3], Section 5] show that $\mathcal{T}_{3,n,\alpha}$ is bijective on the domain $\Omega_{3,n,\alpha}$ up to null sets (and thus our two descriptions of $\Omega_{3,n,\alpha}$ agree). Those arguments discuss the portions of the domain which fiber over the respective cylinders of the interval map, and their images. As Figure 3 hints (in an albeit particularly simple case) these arguments succeed for $m = 2$ and any α interior to matching intervals. This is due, besides the aforementioned success of [[CKS3], Lemma 14], to the fact that those arguments never need to consider the order of C , but are rather based on the geometry of the actions of A and C .

3.7.2. An example: $\Omega_{2,3,1/5}$. Our $\Omega_{2,3,1/5}$, see Figure 3, must agree with the attractor for $\alpha = 1/5$ shown in [[CIT], Figure 8]. That this is so, up to a change of coordinates (due to different normalizing choices), the reader can easily verify from the following detailed information.

With $m = 2, n = 3$, one has $t = t_{2,3} = 1$ and the equality $A^5 C \cdot 1/5 = 0$ is trivially checked. Thus, $\alpha = 1/5$ belongs to the matching interval $J_{k=5,v=1}$. The left endpoint of this matching interval is $\zeta = 3 - 2\sqrt{2} \sim 0.172$, as verified by $A^5 C \cdot \zeta = \zeta - 1$. Since $v = c_1 = 1$, (25) gives $L_{5,1} A = W^4 = (A^{-2} C)^4$. From the definitions in the previous subsection, the upper horizontal boundaries of $\Omega_{2,3,1/5}$ have y -coordinates $\{-\ell_0(\zeta), -\ell_1(\zeta), -\ell_2(\zeta), -\ell_3(\zeta), -\ell_4(\zeta)\} \sim \{0.828, 0.793, 0.739, 0.646, 0.453\}$. Since $\eta = (\sqrt{2} - 1)/2 \sim 0.207$ satisfies $C^{-1} A C A^5 C \cdot \eta = \eta + 1$, this is the right endpoint of the matching interval. Thus the lower horizontal boundaries of $\Omega_{2,3,1/5}$ are $\{-r_0(\eta), -r_1(\eta)\}$ of approximate values $\{-0.207, -0.172\}$.

3.8. Increasing, decreasing, constancy of $\alpha \mapsto \mu(\Omega_{2,n,\alpha})$ on matching intervals. The case of $n = 3$ is treated in [CIT]. We thus assume $n > 3$. We determine the behavior of $\alpha \mapsto \mu(\Omega_{2,n,\alpha})$ in our setting by describing the elementary geometry of the passage from $\Omega_{2,n,\alpha}$ to $\Omega_{2,n,\alpha'}$ for α, α' in the same matching interval. In fact, when α, α' are close neighbors as discussed in [CKS2, CKS3] then a more refined result can be achieved by using the technique of ‘quilting’, see [[CKS2], Main Theorem 2]

Recall that given a matching interval $J_{k,v}$ that $\underline{S}(k, v)$ denotes the length of $\underline{d}(k, v)$ as a word in $\{-1, -2\}$, and $\bar{S}(v)$ is the length of $\bar{d}(k, v)$ as a word in $\{k, k+1\}$.

Lemma 3.4. *Fix $n > 3$. The function $\alpha \mapsto \mu(\Omega_{2,n,\alpha})$ increases on any matching interval $J_{k,v}$ for which $\underline{S}(k, v) < \bar{S}(v)$. The function decreases on $J_{k,v}$ for which the inequality is in the other direction, and is constant when $\underline{S}(k, v) = \bar{S}(v)$.*

Proof. For $\alpha \in J_{k,v}$, the domain $\Omega_{2,n,\alpha}$ is the union of rectangles whose heights are depend only on the maps $T_{2,n,\zeta_{k,v}}$ and $T_{2,n,\eta_{k,v}}$ whereas the widths of these rectangles are the distance between nearest elements of $\{\ell_i(\alpha) \mid 0 \leq i \leq \underline{S}(k, v)\} \cup \{r_j(\alpha) \mid 0 \leq j \leq \bar{S}(v)\}$. Along the $J_{k,v}$ the ordering of this set remains unaltered.

Given $\alpha < \alpha'$ both in $J_{k,v}$ we begin with $\Omega_{2,n,\alpha}$. Let $\mathcal{D} = [\ell_0, \ell'_0] \times [r_{-\bar{S}(v)-1}, y_1]$, with simplified notation of ℓ_0 for $\ell_0(\alpha)$ and ℓ'_0 for $\ell'_0(\alpha)$. Also let $\mathcal{A} = T_A(\mathcal{D})$.

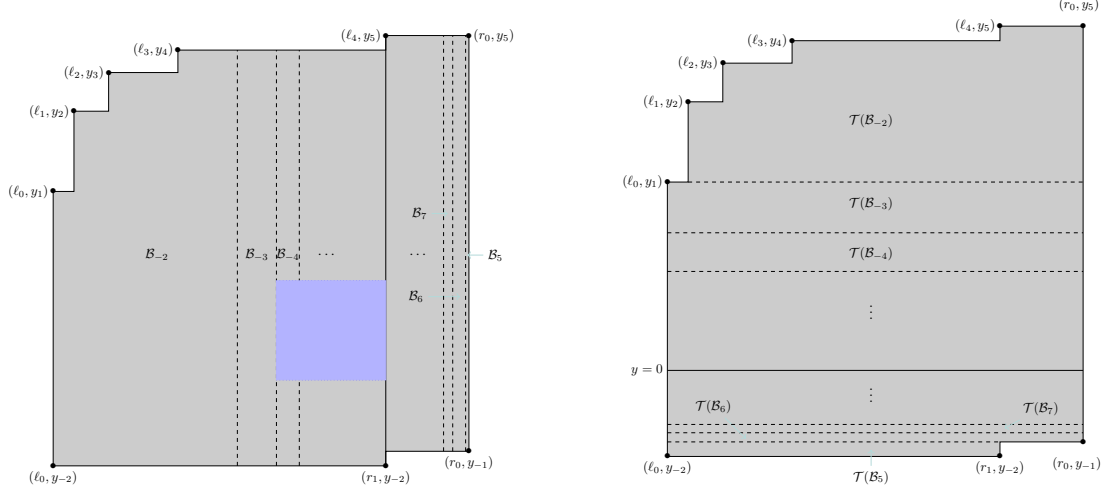


FIGURE 3. On the left: The planar domain $\Omega_{2,3,1/5}$, and its partitioning ‘blocks’; each B_i fibers over the cylinder where $T_{2,3,1/5}$ is given by $x \mapsto A^i C \cdot x$. The $T_{2,3,1/5}$ -orbit of $r_0(1/5) = 1/5$ is labeled by r_0, r_1 ; the $T_{2,3,1/5}$ -orbit of $\ell_0(1/5) = -4/5$ is labeled by $\ell_0, \ell_1, \dots$. One has $r_2(\alpha) = \ell_5(\alpha)$ for other α in the matching interval $J_{k,v} = J_{5,1}$, but here the special equality arises: $r_1 = \ell_4 = 0$. The y -values of the horizontal boundaries are derived from appropriate $T_{2,3,\eta_{5,1}}$ - and $T_{2,3,\eta_{5,1}}$ -orbits, see the discussion of § 3.7.2. The blue rectangle indicates the region $\Omega_{\alpha,d}^+$ of Lemma 3.9. On the right: The image of $\Omega_{2,3,1/5}$ under $T_{2,3,1/5}$ is $\Omega_{2,3,1/5}$ itself.

We claim that

$$\Omega_{2,n,\alpha'} = \left(\Omega_{2,n,\alpha} \setminus \bigcup_{i=0}^{\underline{S}(k,v)} \mathcal{T}_{2,n,\alpha}^i(\mathcal{D}) \right) \cup \bigcup_{i=0}^{\overline{S}(v)} \mathcal{T}_{2,n,\alpha'}^i(\mathcal{A}),$$

see Figure 4. Since $\mathcal{T}_A$ and each application of either $\mathcal{T}_{2,n,\alpha}$ or $\mathcal{T}_{2,n,\alpha'}$ preserves μ -measure, the claim implies $\mu(\Omega_{2,n,\alpha'}) = \mu(\Omega_{2,n,\alpha}) + (\overline{S}(v) - \underline{S}(k,v)) \mu(\mathcal{D})$, whereby the result holds.

We now show that the claim does hold.

Step 1: Correct deletion. Since both $\alpha, \alpha' \in J_{k,v}$, the $T_{2,n,\alpha}$ -digits of $\ell_0(\alpha)$ through $\ell_{\underline{S}(k,v)-1}(\alpha)$ agree with the $T_{2,n,\alpha}$ -digits of $\ell_0(\alpha')$ through $\ell_{\underline{S}(k,v)-1}(\alpha')$ and similarly for the respective digits of $r_0, \dots, r_{\overline{S}(v)-1}$ and $r'_0, \dots, r'_{\overline{S}(v)-1}$. For $0 \leq i \leq \underline{S}(k,v) - 1$ let E_i be the suffix of $L_{k,v}A$ of length i . With this, $\ell_i = T_{2,n,\alpha}^i(\ell_0) = E_i \cdot \ell_0$.

From (26) it follows that for $i \in \{0, \dots, \underline{S} - 1\}$, one has

$$(27) \quad y_{\tau(i+1)} = R A^{p_i} C R^{-1} \cdot y_{\tau(i)},$$

where p_i is the digit of ℓ_i (common throughout $J_{k,v}$), see [[CKS3], Lemma 25 (i)]. It follows that $\mathcal{T}_{2,n,\alpha}^i(\mathcal{D}) = [\ell_i, \ell'_i] \times [e_i, y_{\tau(i)}]$ with $e_i = R E_i R^{-1} \cdot y_{-\overline{S}(v)-1}$. Thus, the union of the $\mathcal{T}_{2,n,\alpha}^i(\mathcal{D})$ accounts exactly for $\Omega_{2,n,\alpha'} \setminus \Omega_{2,n,\alpha}$ if and only if $e_j = y_{\tau(j)-1}$ for all $1 \leq j \leq \underline{S}(k,v)$.

Let p_0 be the initial digit of $\ell_0(\alpha)$, again this is common throughout $J_{k,v}$. That $e_j = y_{\tau(j)-1}$ for all $1 \leq j \leq \underline{S}(k,v)$ is implicit in [[CKS3]], but the proof is not given. We give the proof.

In our case of $n > 3$, [[CKS3], Lemma 33] implies that $\tau(1) - 1 = \tau(1 + i_{\underline{S}})$. From [[CKS3], Lemma 34], $\tau(\underline{S} - i_2) = \underline{S}$ and equivalently $i_{\underline{S}} = \underline{S} - i_2$. Therefore, $\tau(1 + i_{\underline{S}}) = \tau(\underline{S} - (i_2 - 1))$.

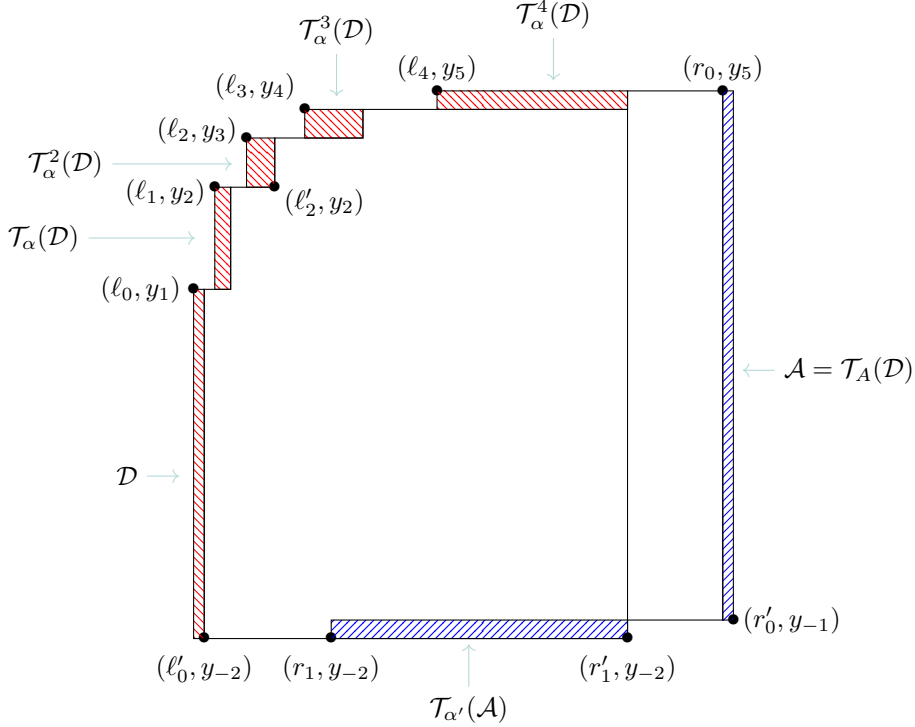


FIGURE 4. **Decreasing mass along $J_{1,5}$.** With $m = 2, n = 3$, both $\alpha = 0.18$ and $\alpha' = 0.2$ lie in the matching interval $J_{1,5}$. Along any matching interval, the coordinates of all tops and bottoms of the rectangles which comprise the various $\Omega_{2,n,\alpha}$ are constant (here $y_1, \dots, y_5$ and y_{-2}, y_{-1}). The function $\mathcal{T}_A$ maps the rectangle $\mathcal{D} = [\ell_0, \ell'_0] \times [r_{-\bar{S}-1}, y_1]$ to $\mathcal{A} = [r_0, r'_0] \times [r_{-1}, y_{\bar{S}+1}]$, where here we use simplified notation of ℓ_0 for $\ell_0(\alpha)$ and ℓ'_0 for $\ell'_0(\alpha)$, and similarly for r_0, r'_0 . To pass from $\Omega_{2,n,\alpha}$ to $\Omega_{2,n,\alpha'}$, we delete the initial $\mathcal{T}_{2,n,\alpha}$ -orbit of $\mathcal{D}$ and add $\mathcal{A}$ and its initial $\mathcal{T}_{2,n,\alpha}$ -orbit. Since each step preserves μ -measure, the sign of the change in mass is that of $(\bar{S}(k) - \underline{S}(k, v))$. See Lemma 7.

By definition, $y_{\tau(\underline{S}-(i_2-1))} = -\ell_{i_2-1}(\zeta_{k,v})$. Thus, $y_{\tau(1)-1} = -\ell_{i_2-1}(\zeta_{k,v})$. Thus (26) implies $(RA^{p_0}CR^{-1})^{-1} \cdot y_{\tau(1)-1} = -A^{p_0}C \cdot \ell_{i_2-1}(\zeta_{k,v})$. By [[CKS3], Corollary 44], $\ell_0(\eta_{k,v}) = \ell_{i_2}(\zeta_{k,v})$ and $\ell_{i_2}(\zeta_{k,v}) = A^{p_0-1}C \cdot \ell_{i_2-1}(\zeta_{k,v})$. Thus, $(RA^{p_0}CR^{-1})^{-1} \cdot y_{\tau(1)-1} = -A \cdot \ell_0(\eta_{k,v}) = -r_0(\eta_{k,v})$. Therefore, $y_{\tau(1)-1} = RA^{p_0}CR^{-1} \cdot -r_0(\eta_{k,v})$.

Now, suppose that some $e_j = y_{\tau(j)-1}$ with $1 \leq j < \underline{S}(k, v) - 1$. By (27), if both $\ell_j(\zeta_{k,v})$ and $\ell_{i_{\tau(j)-1}}(\zeta_{k,v})$ share the same digit, p_j , then $e_{j+1} = y_{\tau(j+1)-1}$. If they do not share the same digit, then by the shape of $\Omega_{2,n,\alpha}$ it can only be that the digit corresponding to $\ell_{i_{\tau(j)-1}}(\zeta_{k,v})$ is rather $p_j - 1$ and in fact $\ell_{i_{\tau(j)}}(\zeta_{k,v})$ is the least entry of the initial orbit segment of $\ell_0(\zeta_{k,v})$ in the p_j cylinder, while $\ell_{i_{\tau(j)-1}}(\zeta_{k,v})$ is the greatest in the $p_j - 1$ cylinder. The first of these is sent by $A^{p_j}C$ to ℓ_2 , the second by $A^{p_j-1}C$ to $\ell_{\underline{S}}$. Thus, $\tau(j+1) = 2$ and $e_{j+1} = RAR^{-1} \cdot y_{\tau(\underline{S})}$. But, [[CKS3], Lemma 34] gives $y_1 = RAR^{-1} \cdot y_{\tau(\underline{S})}$. We thus have that $e_{j+1} = y_{\tau(j+1)-1}$ in this case as well.

Step 2: Correct addition. Since $T_A \cdot x = x + t$, we have $\mathcal{A} = [r_0, r'_0] \times [RAR^{-1} \cdot r_{-\bar{S}(v)-1}, RAR^{-1} \cdot y_1]$. By [[CKS3], Lemma 25 (iv) and Lemma 29 (v)], $\mathcal{A} = [r_0, r'_0] \times [y_{-1}, y_{\bar{S}(k,v)+1}]$. Thereafter, the fact that $-r_0(\zeta_{k,v})$ is the second lowest height, thus that it equals $y_{-\bar{S}}$, follows from [[CKS3], Lemma 45]. This allows arguments directly analogous to those in Step 1 to succeed. $\square$

Definition 3.5. Suppose $n > 3$. The *middle matching interval* is

$$\mathcal{M}_n = \begin{cases} J_{1,(n-2)/2} & \text{if } 2|n; \\ J_{1,(\frac{n-3}{2}, 1, \frac{n-3}{2})} & \text{otherwise.} \end{cases}$$

We show that the middle matching interval does contain the midpoint of the parameter interval, and that the mass is constant along this matching interval. Recall from § 3.3 that for odd n , [BKS] defined R_n to be the positive root of $x^2 - (2 - t_{2,n})x - 1$. If n is even, let $R_n = 1$.

Lemma 3.6. Suppose $n > 3$. The value $\alpha = 1/2$ belongs to $\mathcal{M}_n$. Furthermore,

$$\mathcal{M}_n = \left[1 - \frac{R_n}{t_{2,n}}, 1 + \frac{R_n}{t_{2,n}} \right).$$

Moreover, the function $\alpha \mapsto \mu(\Omega_{2,n,\alpha})$ is constant on $\mathcal{M}_n$, of value

$$\mu(\Omega_{2,n,\alpha}) = \begin{cases} 2 \ln \left[\frac{1 + \cos \pi/n}{\sin \pi/n} \right] & \text{if } 2|n; \\ 2 \ln \left(\frac{1 + R_n}{2 \sin \frac{\pi}{2n}} \right) & \text{otherwise.} \end{cases}$$

Proof. Recall from that for the symmetric Rosen maps, $T_{2,n,1/2}$, we have $T_{2,n,1/2}(-x) = -T_{2,n,1/2}(x)$. As in § 3.2, this allows us to promote results about the Rosen maps themselves to the $T_{2,n,1/2}$. For ease let v_n be defined such that $\mathcal{M}_n = J_{1,v_n}$.

The paper [BKS] shows that R_n gives the highest y -value of their planar domain for the Rosen map of index n . See § 3.2, we then have that this is the highest y -value of $\Omega_{2,n,1/2}$ and also -1 times that value gives the lowest y -value. The basic results of [BKS] and symmetry show that the $T_{2,n,1/2}$ -orbits of $r_0(1/2)$ and $\ell_0(1/2)$ are in agreement with both values being sent to $x = 0$ upon applications of R_{1,v_n} and $L_{1,v_n}A$, respectively. This shows that $\alpha = 1/2$ does lie in $\mathcal{M}_n$, and in fact is the special midpoint value (some author refer to such a value as a ‘pseudocenter’, here it is the center). Symmetry also shows that there is the same number of entries in each of the $T_{2,n,1/2}$ -orbits of $r_0(1/2)$ and $\ell_0(1/2)$ before they meet at $x = 0$.

Since R_n gives the above mentioned highest values, we have that $\ell_0(\zeta_{1,v_n}) = -R_n$. Since for any α one has $\alpha = 1 + \ell_0(\alpha)/t$, the right endpoint of each $\mathcal{M}_n$ is as claimed. Symmetry gives the value of its left endpoint.

Since the $T_{2,n,1/2}$ -orbits of $r_0(1/2)$ and $\ell_0(1/2)$ meet at $x = 0$ after the same number of steps, we have $\bar{S}(v_n) = \underline{S}(1, v_n)$. By Lemma 3.4, $\mu(\Omega_{2,n,\alpha})$ is hence constant along $\mathcal{M}_n$ for each n . As mentioned in § 3.2, $\mu(\Omega_{2,n,1/2})$ equals twice the μ -mass of the planar domain for the Rosen map of index n . With the correction of § 3.3, we find the claimed expressions for $\mu(\Omega_{2,n,\alpha})$. $\square$

Lemma 3.7. Suppose $n > 3$. The function $\alpha \mapsto \mu(\Omega_{2,n,\alpha})$ decreases on any matching interval $J_{k,v}$ such that $J_{k,v}$ lies to the left of $\mathcal{M}_n$, is constant on $\mathcal{M}_n$, and increases on any other matching interval.

Proof. We already have the constancy of mass function $\alpha \mapsto \mu(\Omega_{2,n,\alpha})$ along the middle matching interval, for each n .

For any n , the symmetry $T_{2,n,1-\alpha}(-x) = -T_{2,n,\alpha}(x)$ allows one to find that on each matching interval to the right of $\mathcal{M}_n$ there is a corresponding matching interval to the left of $\mathcal{M}_n$, but

with a reversal of steps in left- and right-orbits (up to matching). In particular, if the mass function $\alpha \mapsto \mu(\Omega_{2,n,\alpha})$ decreases on matching intervals to the left of $\mathcal{M}_n$ then it increases on those to the right of $\mathcal{M}_n$. That is, we need only show the claimed result for matching intervals to the left of $\mathcal{M}_n$.

For any n and any v , from (17) we find that $\bar{S}(v) = \sum_{i=1}^s c_i + \sum_{i=1}^{s-1} d_i$. Each appearance W of (24) in $L_{k,v}$ contributes $n-2$ to $\underline{S}(k,v)$. Thus, each appearance of $\tilde{\mathcal{C}} = \tilde{\mathcal{C}}_k = W^{k-1}(A^{-1}C)^{-1}$ contributes $-1 + (k-1)(n-2)$, and each appearance of $\tilde{\mathcal{D}} = \tilde{\mathcal{C}}_{k+1}$ contributes $-1 + k(n-2)$. From (25),

$$\underline{S}(k,v) = (n-3) + [(k-1)n - 2k + 1](-1 + \sum_{i=1}^s c_i) + [kn - 2k - 1] \sum_{i=1}^{s-1} d_i + (k-1)(n-2).$$

With $n > 3$ and v fixed, this is clearly an increasing function of k . With $k = 2$, we have $\underline{S}(2,v) = (n-3)(\sum_{i=1}^s c_i) + [2n-5] \sum_{i=1}^{s-1} d_i + (n-2)$. Since $n > 3$, this is greater than $\bar{S}(v)$. Thus for all $k \geq 2$ we have that $\underline{S}(k,v) - \bar{S}(v) > 0$. By Lemma 7, we have that the mass function $\alpha \mapsto \mu(\Omega_{2,n,\alpha})$ decreases along any $J_{k,v}$ with $k \geq 2$.

We turn to the setting of $k = 1$. Note that

$$(28) \quad \underline{S}(1,v) - \bar{S}(v) = (n-3)(1 + \sum_{i=1}^{s-1} d_i) + 1 - 2 \sum_{i=1}^s c_i.$$

If $s = 1$ this becomes $\underline{S}(1,c_1) - \bar{S}(v) = n-2-2c_1$. (Recall that $c_1 < n-2$ when $k = 1$.) Thus, $\underline{S}(1,c_1) - \bar{S}(c_1) \geq 0$ if and only if $(n-2)/2 \geq c_1$. Equality holds exactly when $c_1 = (n-2)/2$ and of course this is possible exactly when n is even.

We can now easily show that the case of $n = 4$ holds. Indeed, here $\mathcal{M}_4 = J_{1,1}$, which is a leftmost subinterval of $I_{1,1}$. Hence, any matching interval to the left of $\mathcal{M}_4$ is contained in the union of the $I_{1,k}$, $k \geq 2$. Therefore, we do have that the mass function decreases along all matching intervals to the left of $\mathcal{M}_4$.

For each n , we have that $\underline{S}(1,v) - \bar{S}(v)$ is zero when v is the defining word of $\mathcal{M}_n$. As well, for fixed v , (28) shows that $\underline{S}(1,v) - \bar{S}(v)$ strictly increases with n . Thus for $n = 5$ we must merely show that $\underline{S}(1,v) - \bar{S}(v)$ is positive between $J_{1,1}$ and $\mathcal{M}_5 = J_{1,111}$. But, as [CKS] shows, the matching intervals between these two are the $J_{1,1q1}$ with $q \geq 2$ and the $J_{1,v}$ with v running through the descendants of these various $1q1$. Now, with $n > 4$ we have that $\underline{S}(1,1) - \bar{S}(1)$ is positive, and to pass from $v = 1$ to $v = 1q1$, we add q to the sum of the d_j and 2 to that of the c_j . Again from (28), we find $(\underline{S}(1,1q1) - \bar{S}(1q1)) - (\underline{S}(1,1) - \bar{S}(1)) = (n-3)q - 4$. Since $n = 5$, with $q \geq 2$ this is positive. In any $v \in \mathcal{V}$ with $c_1 = 1$ all of the $c_j = 1$. A consequence of the construction of the tree of words $\mathcal{V}$, see Definition 2.5, is that upon fixing such a q , any descendant of $1q1$ has its d_j all in $\{q, q+1\}$. Since any descendant is formed by appending a string of alternating d_j and c_{j+1} , we see that the difference similar to the above is again positive. That is, the result also holds for $n = 5$.

Now suppose that the result holds for some odd $n = 2h + 3$ with $h \geq 1$ and consider $n' = n + 1 = 2(h+2)$. We aim to show that, with n' as our indexing value, $\underline{S}(1,v) - \bar{S}(v) > 0$ for all matching intervals between $J_{1,h1h}$ and $J_{1,h+1}$, knowing that $\underline{S}(1,h1h) - \bar{S}(h1h)$ is positive. These intermediate matching intervals are all of the form $J_{1,v}$ with v a descendant of $h1h$.

If $h = 1$ then as above each such descendant v has all $c_j = 1$ and all $d_j \in \{1, 2\}$ we pass from 111 to v by appending an equal number of c_j and d_j values. The positivity of $\underline{S}(1,v) - \bar{S}(v)$ thus holds. If $h \geq 2$ then each descendant has all $d_j = 1$ and all $c_j \in \{h, h+1\}$. Thus, the change in the difference of the right hand side of (28) as we add one pair of d_j, c_{j+1} at a time

is $(n' - 3)d_j - 2c_{j+1} \geq n' - 3 - 2(h + 1) = 2(h + 2) - 3 - 2(h + 1) = 0$. We conclude that the difference $\underline{S}(1, v) - \overline{S}(v)$ is positive for all of these v .

We now aim to show the result for $n'' = n + 2 = 2(h + 1) + 3$. Hence, we must show, with our value n'' , the positivity along the matching intervals between $J_{1, h+1}$ and $J_{1, (h+1)1(h+1)}$, knowing that (28) is positive when $v = h + 1$. That is, we must show positivity when v is any descendant of $h + 1$ that is not a descendant of $(h + 1)1(h + 1)$. The children of $h + 1$ are of the form $(h + 1)(1h)^q 1(h + 1)$. Thus for each $q > 0$ we add $(q + 1)h + 1$ in total to the sum of the c_j -values, and $q + 1$ to the sum of the d_j . Thus (28) changes by $(n'' - 3)(q + 1) - 2(q + 1)h - 2 = q - 1$. Any descendant of $v = (h + 1)(1h)^q 1(h + 1)$ with $q > 0$ can be formed by repeatedly appending pairs of the form $1h$ or $1(h + 1)$ to v . Each such step results in the difference changing by $(n'' - 3) - 2h = 2$ or by 1. Thus, the result holds for n'' .

Thus, an induction proof shows the result holds for all $n > 3$. $\square$

3.9. Continuity of $\alpha \mapsto \mu(\Omega_{2,n,\alpha})$. The table [[CKS3], Table 1] also points to definitions of $\Omega_{3,n,\alpha}$ in the cases that α is either endpoint of a matching interval; in short, when compared to $\Omega_{3,n,\alpha}$ with α in the interior of the matching interval, at the left endpoint there is a ‘missing’ leftmost lower rectangle, and at the right a missing rightmost upper rectangle, see [[CKS3], Figures 7, 8]. These definitions are easily adjusted to the case of $m = 2$, and the proofs given in [CKS3] succeed.

The remaining values of α are non-matching values. The arguments of [[CKS3], § 6] can be directly adapted to the setting of $m = 2$. There, one has that $\Omega_{3,n,\alpha}$ includes a central rectangle see [[CKS3], Figure 9], given by the union of the full cylinders $\Delta_{3,n,\alpha}(k')$ times a constant y -fiber Φ equal to the closure of $\mathbb{I}_{3,n,\alpha}$. The full $\Omega_{3,n,\alpha}$ is the closure of the union of the $\mathcal{T}_{3,n,\alpha}$ -orbit of a certain $\mathcal{Z}$ formed by the union of the central rectangle with the left remaining subinterval times a constant y -fiber Φ^- equal to negative one times the closure of the right remaining subinterval and similarly for a remaining rectangle of y -fiber denoted Φ^+ . Each non-matching value is the limit of endpoints of matching intervals; for each such endpoint α' , there is an analogous $\mathcal{Z}_{\alpha'}$. The various $\mathcal{Z}_{\alpha'}$ nicely converge to $\mathcal{Z}_\alpha$ and one shows that the union of the images under the α' -map does indeed converge to the union of the $\mathcal{T}_{3,n,\alpha}$ -orbit of $\mathcal{Z}_\alpha$. All of this, in adapted form, carries over to show the continuity of $\alpha \mapsto \mu(\Omega_{2,n,\alpha})$ when $0 < \alpha < 1$.

Corollary 7. *Fix $n \geq 3$. The function $\alpha \mapsto \mu(\Omega_{2,n,\alpha})$: decreases on $0 < \alpha < 1 - \frac{R_n}{t_{2,n}}$; is constant on $[1 - \frac{R_n}{t_{2,n}}, 1 + \frac{R_n}{t_{2,n}})$; and, increases on $1 + \frac{R_n}{t_{2,n}} \leq \alpha < 1$. Its maximal value is as given in Lemma 3.6.*

Proof. Since the mass function is continuous on $(0, 1)$, the result follows from Lemmas 3.6 and 3.7. $\square$

Since the maximal value of $\mu(\Omega_{2,n,\alpha})$ goes to zero as $n \rightarrow \infty$, we have completed the proof of Theorem 2.

3.10. Planar extensions for $\alpha = 0, 1$. We also give planar extensions for the endpoints of the parameter intervals.

3.10.1. Planar extensions $\Omega_{2,n,\alpha}$ for $\alpha = 0$. For each $n \geq 3$, the domain $\Omega_{3,n,0}$ is given in [CS]. We now sketch the adjustment to the setting of $m = 2$.

Fix $m = 2$ and $n \geq 3$, and set $t = t_n = t_{2,n} = 2\cos(\pi/n)$. Recycle notation and let $g(x) = kt - 1/x$, with k the unique integer such that $g(x) \in [-t, 0)$. Now, $-t \leq kt - 1/x < 0$ implies $-(k + 1)t \leq -1/x < -kt$ (and since $x < 0$ we have $-1/x > 0$ so $k < 0$). That is, $\frac{1}{(k+1)t} \geq x > \frac{1}{kt}$. We find that the cylinders corresponding to $-k, k \geq 2$ are all full, whereas (if $n > 3$) the cylinder of $k = -1$ is sent to $[-t + 1/t, 0)$. One finds that the orbit of $x = -t$ is

periodic of length $n-2$ and corresponds to following the admissible word $W = A^{-2}C(A^{-1}C)^{n-3}$. The reader may be reassured by the facts that both $[[\text{CKS}], \text{equation (2.1)}]$ holds when $m = 2$, and (projectively) $W = \begin{pmatrix} 1+t^2 & t^3 \\ -t & 1-t^2 \end{pmatrix}$ as in $[[\text{CS}], \text{equation (3.1)}]$.

It is easily verified that $\mathcal{T}_{2,n=3,0}$ is bijective up to null sets on $\Omega_{2,3,0} := [-1, 0] \times [0, 1]$. We thus assume $n > 3$ in what follows. Let $\ell_i = g^i(-t)$ for $0 \leq i \leq n-3$. We have $-t = \ell_0 < \ell_1 < \dots < \ell_{n-3}$ and $\{\ell_0, \dots, \ell_{n-4}\} \subset \Delta(-1)$ while $\ell_{n-3} \in \Delta(-2)$, where we use simplified notation for cylinders. Since $A^{-2}C \cdot \ell_{n-3} = -t$, in fact $\ell_{n-3} = C \cdot t = -1/t$.

Now let $y_1 = 1/t$, a value fixed by RWR^{-1} . Further, let $y_i = RA^{-1} \cdot y_{i-1}$ for $1 \leq i \leq n-2$. We find that $RA^{-2} \cdot y_{n-2} = y_1$. Thus, $y_{n-2} = t$. Let

$$\Omega = \Omega_{2,n,0} = ([\ell_{n-3}, 0] \times [0, t]) \cup \bigcup_{i=0}^{n-4} ([\ell_i, \ell_{i+1}] \times [0, y_{i+1}]),$$

see Figure 5 for the case of $n = 5$.

We have that $\mathcal{S}' : (x, y) \mapsto (-y, -x)$ sends the pair (ℓ_{n-3}, y_{n-2}) to (ℓ_0, y_1) . It is straightforward to check that for any invertible $M = \begin{pmatrix} a & b \\ -b & d \end{pmatrix}$ that $\mathcal{T}_M(\mathcal{S}'(x, y)) = \mathcal{S}' \circ \mathcal{T}_{M^{-1}}(x, y)$ for any (x, y) . Since $(\ell_1, y_2) = \mathcal{T}_{A^{-1}C}(\ell_0, y_1)$ and $(\ell_{n-4}, y_{n-3}) = \mathcal{T}_{(A^{-1}R)^{-1}}(\ell_{n-3}, y_{n-2})$, by recursion we find that the vertices of Ω are preserved by $\mathcal{S}'$. That is, Ω is preserved by $\mathcal{S}'$. Notice that when n is odd we have $(\ell_{(n-3)/2}, y_{(n-1)/2}) = (-1, 1)$. Since there are vertices lying along the curve $y = -1/x$, one sees that $\mu(\Omega)$ is infinite.

3.10.2. Planar extensions for accelerated $T_{2,n,0}$. Similar to $[\text{CS}]$, we now accelerate $g = T_{2,n,0}$ by inducing past the cylinder of W . We treat only the case of $n > 3$ and hence $t > 1$, the remaining case is readily handled. Let $\epsilon_0 = W^{-1} \cdot 0 = t^3/(1-t^2)$ and define $f : [-t, 0) \rightarrow [-t, 0)$ by $f(x) = g(x)$ when $x > \epsilon_0$ and for $-t < x \leq \epsilon_0$ let $f(x) = W^{j(x)} \cdot x$ where $j(x)$ is minimal such that $W^{j(x)} \cdot x > \epsilon_0$.

Let $\mathcal{F}$ be the planar map associated, in our usual manner, to f . We give a bijectivity domain for $\mathcal{F}$. For this, let ϵ_i with $1 \leq i \leq n-3$ be the initial g -orbit of ϵ_0 , that is, $\epsilon_i = (A^{-1}C)^i \cdot \epsilon_0$. Since $W \cdot \epsilon_0 = 0$, we have that ϵ_{n-3} is the right endpoint of the cylinder $\Delta(-2)$. Also, $\ell_i < \epsilon_i < \ell_{i+1}$ for $0 \leq i \leq n-4$ and $\ell_{n-3} < \epsilon_{n-3}$; recall that ℓ_{n-3} is sent to $-t$ and hence it is the left endpoint of $\Delta(-2)$. Let $z_1 = RWR^{-1} \cdot 0$. Note that $\mathcal{T}_{W^{k+1}}$ sends the rectangle $(W^{-k-1} \cdot \epsilon_0, W^{-k} \cdot \epsilon_0) \times [0, z_1]$ to lie directly above the $\mathcal{T}_{W^k}$ -image of $(W^{-k} \cdot \epsilon_0, W^{-k-1} \cdot \epsilon_0)$. Since y_1 is a fixed point of RWR^{-1} one finds the union over all k of these images is the rectangle $[\epsilon_0, 0] \times [z_1, y_1]$. With this as a guide, we now let

$$\Omega = ([\ell_0, \epsilon_0] \times [0, z_1]) \cup ([\epsilon_{n-3}, 0] \times [0, t]) \cup \bigcup_{i=0}^{n-4} ([\epsilon_i, \epsilon_{i+1}] \times [0, y_i]),$$

see Figure 5 for the case of $n = 5$.

One finds that $\Omega = \Omega_{2,n,0} \setminus (\mathcal{D}_1 \cup \mathcal{D}_2)$ where

$$(29) \quad \mathcal{D}_1 = \bigcup_{i=0}^{n-3} \mathcal{T}_{2,n,0}^i([\ell_0, \epsilon_0] \times [z_1, y_1]) \text{ and } \mathcal{D}_2 = \bigcup_{i=1}^{n-3} \mathcal{T}_{2,n,0}^i([\ell_0, \epsilon_0] \times [0, z_1]).$$

Given that $\mathcal{F}$ maps $[\ell_0, \epsilon_0] \times [0, z_1]$ bijectively to $[\epsilon_0, 0] \times [z_1, y_1]$ and the bijectivity of $\mathcal{T}_{2,n,0}$ on $\Omega_{2,n,0}$, the bijectivity of $\mathcal{F}$ on Ω will follow as soon as we show that $RA^{-2}CR^{-1}$ maps y_{n-3} to z_1 . Given the matrix form of W , we find $z_1 = RW \cdot \infty = -1/((1+t^2)/-t) = t/(1+t^2)$. On the other hand, $y_{n-3} = (RA^{-1}CR^{-1})^{-1} \cdot t$ and thus $RA^{-2}CR^{-1} \cdot y_{n-3} = RA^{-1}R^{-1} \cdot t = -1/((-t-1/t))$. The equality holds and hence Ω is a domain of bijectivity for $\mathcal{F}$.

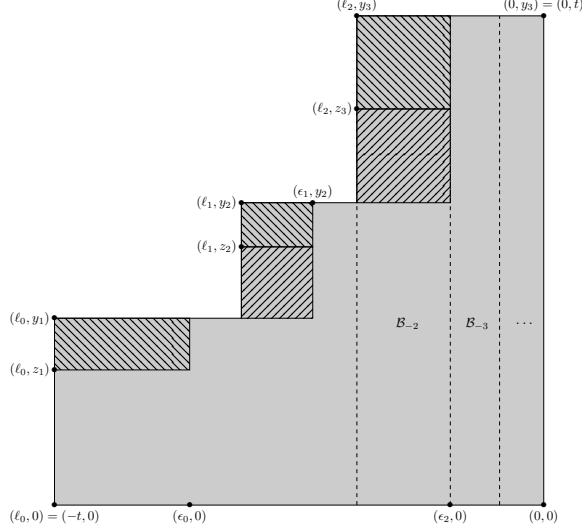


FIGURE 5. The domain $\Omega_{2,5,0}$ is the union of the gray and the hatched regions. These latter are excised to leave the planar domain of the two-dimensional system defined by the interval map given by by accelerating $T_{2,5,0}$ by applying appropriate powers of W for x in $(\ell_0, \epsilon_0) = (-t_{2,5}, \epsilon_0)$, where $\epsilon_0 = W^{-1} \cdot 0$. The negative sloped hatched rectangles form $\mathcal{D}_1$ of (29), the positive sloped hatched rectangles form $\mathcal{D}_2$.

3.10.3. *Planar extension for $(m, n, \alpha) = (2, n, 1)$.* Applying the map $\mathcal{S}$ to $\Omega_{2,n,0}$ results in $\Omega_{2,n,1}$ and thus one can verify the following.

In [[CKS], Section 3] one finds a discussion of $\Omega_{3,n,1}$ for $n \geq 3$. For our setting of $m = 2$, we must check some of that argument, since here $\mu = 0$, but $1/\mu$ appears in that earlier discussion. Recycling notation, we have $\mathbb{I} = [0, t)$ and the non-full cylinder $\Delta(1) = [1/t, t)$ whose image under f is $AC \cdot [1/t, t) = t + [-t, -1/t) = [0, t - 1/t) = [0, (t^2 - 1)/t)$. The other cylinders are full; on the cylinder index by $k > 1$, T takes the form $x \mapsto kt - 1/x$.

The argument of [[CKS], Section 3] goes through to show that $A^2C(AC)^{n-3} = A(AC)^{n-2} = A(CA^{-1})^2 = ACA^{-1}CA^{-1}$ fixes t and the orbit of t is given by $t \mapsto AC \cdot t \mapsto (AC)^2 \cdot t \mapsto \dots \mapsto (AC)^{n-3} \cdot t$, followed by $A^2C \cdot ((AC)^{n-3} \cdot t) = t$. We thus let $s_i = (AC)^i \cdot t$, $0 \leq i \leq n-3$ and set

$$(30) \quad \Omega = \bigcup_{i=1}^{n-2} ([s_i, s_{i-1}] \times [-1/s_{i-1}, 0]) = ([0, 1/t) \times [-t, 0]) \cup \bigcup_{i=1}^{n-3} ([s_i, s_{i-1}] \times [-1/s_{i-1}, 0]).$$

For $k > 1$, by [[CKS2], Proposition 12], the two dimensional $\mathcal{T}$ maps the rectangle $\Delta(k) \times [-t, 0]$ below the image of $\Delta(k+1) \times [-t, 0]$ so that they share a horizontal line. Since $RA^kCR \cdot y = RA^k \cdot y = -1/(kt+y)$ we find that the union of these images, with $2 \leq k < \infty$ is $\mathbb{I} \times [-1/t, 0)$. One easily verifies that $\mathcal{T}_{AC}$ sends $[s_i, s_{i-1}] \times [-1/s_{i-1}, 0]$ to $[s_{i+1}, s_i] \times [-1/s_i, -1/t]$ for $i < n-2$. Hence, Ω is a bijectivity domain for $\mathcal{T}$ up to null sets. We thus hereafter denote it by $\Omega_{2,n,1}$.

Note that $\Omega_{2,3,1}$ is the square $[0, 1] \times [-1, 0]$.

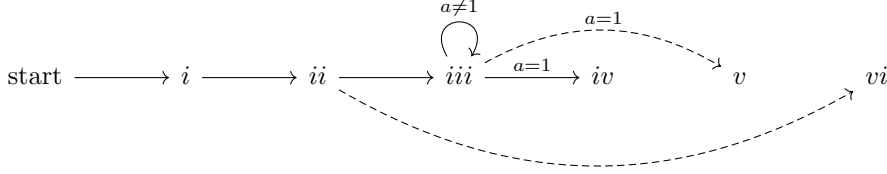


FIGURE 6. Sequence of word processing relations to pass from T_α -admissible word to T_0 -admissible word, used in proof of Lemma 3.8. Processing halts when the demarcation symbol is to the left of the word.

3.11. Continuity of entropy, constancy of ‘Rohlin integrals’.

3.11.1. *Agreement of first return maps.* Using ‘word processing’ we show that for each $0 \leq \alpha < 1$ that upon an appropriate subset of negative $x \in \mathbb{I}_\alpha$ first returns under T_α are given by compositional powers of T_0 . Compare the following with [[CKS2], Lemma 142].

Lemma 3.8. *Fix $m = 2$ and $n \geq 3$. For each $\alpha \in [0, 1)$ and suppose $x \in \mathbb{I}_\alpha$ is negative and that there exists $m \geq 0$ such that both $T_\alpha^m(x) < 0$ and $T_\alpha^{m+1}(x) < 0$. Then there is some $k \in \mathbb{N}$ such that $T_0^k(x) = T_\alpha^{m+1}(x)$. Furthermore, for any y one has $\mathcal{T}_0^k(x, y) = \mathcal{T}_\alpha^{m+1}(x, y)$.*

Proof. With minor adjustments, the proof of [[CKS2], Lemma 142] succeeds in our setting, as we now show.

We can and do assume that m is minimal with respect to the hypotheses. From the definition of our maps, we can write

$$T_{\alpha}^{m+1}(x) = A^{a_{m+1}}CA^{a_m}C \cdots A^{a_1}C \cdot x,$$

with both a_{m+1}, a_1 negative. If $m = 0$, then $A^{a_1}C^{c_1} \cdot x = T_0(x)$ and $k = 1$. We now assume that $m > 0$ and perform word processing on $A^{a_{m+1}}CA^{a_m}C \dots A^{a_1}C$ with the goal to achieve an expression for this element of the group G_n that is in an admissible form for T_0 .

We process from right to left, using the following substitution rules. We start with the augmented word ' $A^{a_{m+1}}CA^{a_m}C\cdots A^{a_1}C$;' and use the substitutions:

$$\begin{array}{llll}
(i) & A^{a_1}C; & \mapsto & A; A^{a_1-1}C \\
(ii) & A^aCA; & \mapsto & ACA; W^{a-1} \\
(iii) & A^aCAC A; & \mapsto & \begin{cases} ACA; W^{a-2}A^{-2}C(A^{-1}C)^{n-4} & \text{if } a \geq 2, n \geq 4; \\ ACA; W^{a-2}A^{-1} & \text{if } a \geq 2, n = 3; \\ ; A^{-2}C(A^{-1}C)^{n-4} & \text{if } a = 1. \end{cases} \\
(iv) & A^aC; A^bC & \mapsto & ACA; W^{a-1}A^{b-1}C \quad \text{if } a > 0 \\
(v) & A^{a_{m+1}}C; & \mapsto & ; A^{a_{m+1}}C \\
(vi) & A^{a_{m+1}}CAC A; & \mapsto & ; A^{a_{m+1}-1}C(A^{-1}C)^{n-3}.
\end{array}$$

Each rule, up to ignoring the *marker* “;”, is the result of applying a group identity. Mainly, it is a version of $A^k CA = CAW^k$ from Lemma 3.3. After a rule is applied, we multiply on the left by the next unprocessed $A^{a_i}C$, and then determine which rule is then to be applied. (In the ‘middle case’ of (iii), we also first multiply the A^{-1} with the term to its right.) The processing halts upon applying (v) or (vi). One easily verifies that processing does halt, confer Figure 6. The result is an admissible T_0 -word, which acts so as to send x to $T_\alpha^{m+1} \cdot x$.

□

We follow the next step in the arguments of [CKS2], and seek a region for which the *first* return under $\mathcal{T}_0$ agrees with that of $\mathcal{T}_\alpha$. To this end, we restrict the region to which we return, so that $\{-1, -2\}$ digits introduced in the word processing could never index an ‘early return’ of a $\mathcal{T}_0$ -orbit.

Lemma 3.9. *Suppose that $0 < \alpha < 1$. Let $d = \min\{-4, -2 + d_\alpha(\ell_0(\alpha))\}$, $N_d = \bigcup_{i \leq d} \Delta_\alpha(i)$ and $\Omega_{\alpha,d}^+ = \Omega_\alpha \cap \{(x, y) \mid y > 0, x < 0, x \in N_d \text{ and } x' \in N_d, \text{ where } \mathcal{T}_\alpha^{-1}(x, y) = (x', y')\}$. The first return maps of $\mathcal{T}_0$ and of $\mathcal{T}_\alpha$ to $\Omega_{\alpha,d}^+$ agree.*

Proof. Since $\Omega_{\alpha,d}^+$ lies in the second quadrant of the (x, y) -plane, one finds $\Omega_{\alpha,d}^+ \subset \Omega_0$ and hence both return maps are defined. Suppose that $(x, y) \in \Omega_{\alpha,d}^+$ and $\mathcal{T}_\alpha^{m+1}(x, y)$ is its first $\mathcal{T}_\alpha$ -return. From the previous lemma, this first return is also *some* return under $\mathcal{T}_0$.

If $m = 0$ then certainly $T_\alpha(x) = T_0(x)$, and the first return of the two-dimensional maps also agree. Now assume that $m > 0$. Suppose

$$T_\alpha^{m+1}(x) = A^{a_{m+1}} C^{c_{m+1}} A^{a_m} C^{c_m} \dots A^{a_1} C^{c_1} \cdot x.$$

Of course, if all a_i are negative, then the orbits are the same and the result clearly holds. In general, there may be intermediate consecutive appearances of negative exponents. We thus add:

$$(vii) \quad A^{a'} C A^a C; \mapsto A^{a'} C; A^a C \text{ if } a < 0, a' < 0,$$

and process so as to achieve an admissible T_0 expression.

Whenever any of our substitution rules inserts T_0 -steps, these insertions are always of corresponding 0-digits in $\{-1, -2\}$. Of course, any x -value with either of these as its 0-digit cannot have α -digit less than or equal to d . Thus, these insertions cannot cause an early $\mathcal{T}_0$ return.

Our rules also make changes in the exponents of appearances of A . However, any such exponent is changed at most once throughout the process, and if it is changed then it is decreased exactly by one. Since each substitution realizes a group identity, it follows that any element in N_d that is in the T_0 -orbit segment of x but not already present in the T_α -orbit segment is either (1) isolated (with respect to this property), or (2) is sent by T_0 to $T_\alpha^m(x)$ while its T_0 -orbit predecessor is not in N_d . But, by hypothesis, $\mathcal{T}_\alpha^m(x, y) \notin \Omega_{\alpha,d}^+$. Therefore the result holds. □

3.11.2. Entropy times mass is constant, proof of Theorem 2.

Proposition 3.10. *Fix $n \geq 3$. The function $\alpha \mapsto h(T_{2,n,\alpha})\mu(\Omega_{2,n,\alpha})$ is constant on $(0, 1)$.*

Proof. Fix $\alpha, \alpha' \in (0, 1)$ and let $\Omega_{\alpha,\alpha'} = \Omega_{\alpha,d}^+ \cap \Omega_{\alpha',d'}^+$, with the natural interpretation of the notation of Lemma 3.9. That lemma implies that the first returns to this region of $\mathcal{T}_{2,n,\alpha}$ and $\mathcal{T}_{2,n,\alpha'}$ agree. These hence define the same dynamical system, and Abramov’s formula yields that their entropy can be expressed as

$$\frac{h(\mathcal{T}_{2,n,\alpha})\mu(\Omega_{2,n,\alpha})}{\mu(\Omega_{\alpha,\alpha'})} = \frac{h(\mathcal{T}_{2,n,\alpha'})\mu(\Omega_{2,n,\alpha'})}{\mu(\Omega_{\alpha,\alpha'})}.$$

Since the entropy of a dynamical system equals the entropy of the natural extension system, the result holds. □

Theorem 2 now follows from Proposition 3.10 and Corollary 7.

3.11.3. *Rohlin integrals are constant.* We continue to follow [CKS3]. Recall the cocycle τ from Definition 2.1.

Definition 3.11. Fix $n \geq 3$. For $\alpha \in [0, 1]$ and for each $x \in \mathbb{I}_{2,n,\alpha}$, let $\tau_\alpha(x) = \tau_{2,n,\alpha}(x) = -2 \log |cx + d|$ where $T_{2,n,\alpha}(x) = (ax + b)/(cx + d)$. Of course, this is simply $\tau_\alpha(x) = \tau(M, x)$ where $T_{2,n,\alpha}(x) = M \cdot x$.

For any n, α , we call the integral $\int_{\Omega_{2,n,\alpha}} \tau_{2,n,\alpha}(x) d\mu$ a *Rohlin integral*. See the first displayed equation in the following proof for a justification of this term.

We now show, in terms of simplified notation, that the integral of $\tau_\alpha(x)$ over Ω_α with respect to $d\mu$ gives the product of $h(T_\alpha)$ with $\mu(\Omega_\alpha)$. Furthermore, these integrals are constant with respect to α , including at the endpoints, this although both $\mu(\Omega_0)$ and $\mu(\Omega_1)$ are infinite.

Lemma 3.12. Fix $n \geq 3$. For $0 < \alpha < 1$,

$$h(T_{2,n,\alpha}) \mu(\Omega_{2,n,\alpha}) = \int_{\Omega_{2,n,\alpha}} \tau_{2,n,\alpha}(x) d\mu = \int_{\Omega_{2,n,0}} \tau_{2,n,0}(x) d\mu = \int_{\Omega_{2,n,1}} \tau_{2,n,1}(x) d\mu.$$

Proof. For $0 < \alpha < 1$, Rohlin's formula gives

$$h(T_{2,n,\alpha}) = \int_{\mathbb{I}_{2,n,\alpha}} \log |T'_{2,n,\alpha}(x)| d\nu_{2,n,\alpha} = \int_{\mathbb{I}_{2,n,\alpha}} \tau_{2,n,\alpha}(x) d\nu_\alpha = \frac{\int_{\Omega_{2,n,\alpha}} \tau_{2,n,\alpha}(x) d\mu}{\mu(\Omega_{2,n,\alpha})},$$

where the last equality holds because $\nu_{2,n,\alpha}$ is the marginal measure for the probability measure on $\Omega_{2,n,\alpha}$ induced by μ . Since $\mu(\Omega_{2,n,\alpha}) < \infty$, our first equality holds.

By Theorem 4 and the discussion of symmetry in § 3.1, it suffices to prove only the first of the two remaining equalities.

We employ a tower construction, with base $\Omega_{\alpha,d}^+$. Throughout, we use the convention that equalities of sets are all considered up to μ -null sets. Poincaré recurrence applied to each of $\mathcal{T}_{2,n,\alpha}^{\pm 1}$ on this positive measure subspace of $\Omega_{2,n,\alpha}$ shows the bijectivity of the first return of $\mathcal{T}_{2,n,\alpha}$ to $\Omega_{\alpha,d}^+$. For each $k \in \mathbb{N}$ we define $R_k \subset \Omega_{\alpha,d}^+$ as the set of points whose first return to $\Omega_{\alpha,d}^+$ is given by applying $\mathcal{T}_{2,n,\alpha}^k$. The R_k thus partition $\Omega_{\alpha,d}^+$ and also $\Omega_{\alpha,d}^+ = \cup_{k \in \mathbb{N}} \mathcal{T}_{2,n,\alpha}^k(R_k)$ up to μ -null sets. Now set $S = \cup_{k \in \mathbb{N}} \sqcup_{i=0}^{k-1} \mathcal{T}_{2,n,\alpha}^i(R_k)$. We have that $\mathcal{T}_{2,n,\alpha}$ bijectively maps the tower S to itself. By the ergodicity of $\mathcal{T}_{2,n,\alpha}$, we deduce that $S = \Omega_{2,n,\alpha}$.

Let ρ denote projection onto the x -coordinate of the first return map on $\Omega_{\alpha,d}^+$. Due to the cocycle property of $\tau(M, x)$ and the additivity of integration, one deduces that

$$\int_{\Omega_{2,n,\alpha}} \tau_{2,n,\alpha}(x) d\mu = \int_{\Omega_{\alpha,d}^+} \log |\rho'(x)| d\mu.$$

Since $\Omega_{2,n,0}$ has infinite measure, we cannot simply repeat this argument to directly achieve the equality with $\Omega_{2,n,0}$ replacing $\Omega_{2,n,\alpha}$. Rather, we turn to the analog of the induced system that [CS] associates to the $\mathcal{T}_{3,n,0}$. Associated to the accelerated system given in § 3.10.1 is the induced system on $\mathcal{I} = \Omega_{2,n,0} \setminus \cup_{i=0}^{n-3} \mathcal{T}_{2,n,0}^i(\mathcal{D})$, where $\mathcal{D}$ is the domain fibered over $[-t, \epsilon)$. Since $\Omega_{2,n,0}$ meets the curve $y = -1/x$ exactly on the initial $n - 2$ points of the $\mathcal{T}_{2,n,0}$ -orbit of $(-t, 1/t)$, this induced system is a finite μ -measure dynamical system. Now, $\mathcal{D}$ is contained in the $\mathcal{T}_{2,n,0}$ -image of the blocks of digit -3 and less. That is, the forward $\mathcal{T}_0$ orbits of the points of $\mathcal{I}$ up to first return to $\mathcal{I}$ gives all of $\Omega_{2,n,0}$. (As usual, these statements are true up to null sets, for example the purely periodic orbit of $(-t, 1/t)$ is not contained in those forward orbits.) But, again by the cocycle property of $\tau(M, x)$ and the additivity of integration, the integration of the logarithm of the absolute value of the first derivative of the accelerated one-dimensional map over $\mathcal{I}$ equals $\int_{\Omega_0} \tau_{2,n,0}(x) d\mu$.

Now, for our $0 < \alpha < 1$, due to the restriction on the digits in the definition of the region $\Omega_{\alpha,d}^+$, the region is contained in $\mathcal{I}$. Hence, Lemma 3.9 implies that the first return maps of these two finite area systems agree. Again the cocycle property gives that the Rohlin integrals of these systems are equal. By transitivity of equality, we have $\int_{\Omega_{2,n,\alpha}} \tau_{2,n,\alpha}(x) d\mu = \int_{\Omega_{2,n,0}} \tau_{2,n,0}(x) d\mu$. $\square$

Lemma 3.12 combines with (22) to achieve the proof of Theorem 5.

4. FROM $\Omega_{2,n,1}$ TO $\Omega_{3,n,1}$

4.1. The domain of bijectivity $\Omega_{3,n,1}$. In [[CKS2], Proposition 12] one finds that for each $n \geq 3$ the region which in our notation is $\Omega_{3,n,1}$ is

$$(31) \quad \Omega_{3,n,1} = ([0, 1] \times [-1, 0]) \cup \bigcup_{i=1}^{n-2} [r_i, r_{i-1}] \times [-1/r_{i-1}, 0],$$

where $r_0, r_1, \dots, r_{n-2}$ is the $T_{3,n,1}$ -orbit of $r_0 = t = t_{3,n}$. See [[CKS2], Figure 3].

Directly above their proposition, they state what we can express, with C_3 dependent on $m = 3$, and A_3 denoting the translation by $t_{3,n}$, as $r_i = (A_3 C_3^2)^{n-2} \cdot t$ for $1 \leq i \leq n-2$, with $r_{n-2} \cdot t = 1$; and, $T^{n-1}(t) = A_3 C_3 (A_3 C_3^2)^{n-2} \cdot t = t$.

4.2. Sending $\Omega_{2,n,1}$ to the right hand part of $\Omega_{3,n,1}$. Fix $n \geq 3$, and note that $t_{3,n} = 1 + t_{2,n}$.

Fix $M_1 = \begin{pmatrix} 1 & 1 \\ 0 & 1 \end{pmatrix}$. We now consider the $\mathcal{T}_M$ image of $\Omega_{2,n,1}$. This function acts on x -coordinates as translation by 1. We have $RM_1 R \cdot y = -1/(1 - 1/y)$. Thus, $\mathcal{T}_{M_1}(t_{2,n}, -1/t_{2,n}) = (t_{3,n}, -1/t_{3,n})$.

Now, $C_3 = \begin{pmatrix} 1 & -1 \\ 1 & 0 \end{pmatrix}$, and thus $C_3^2 M_1 = C_2$, where of course we are letting $C_2 = \begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix}$.

Since $A_3 = M_1 A_2$ (where we let A_2 denote that translation by $t_{2,n}$), we have that $M_1^{-1} A_3 C_3^2 M_1 = A_2 C_2$. Thus, $\mathcal{T}_{M_1}$ conjugates $\mathcal{T}_{A_2 C_2}$ to $\mathcal{T}_{A_3 C_3^2}$. Referring back to (30) and its discussion, we find that $\mathcal{T}_{M_1}$ sends $\Omega_{2,n,1}$ to $\Omega_{3,n,1} \setminus ([0, 1] \times [-1, 0])$. In other words, $\Omega_{3,n,1} = ([0, 1] \times [-1, 0]) \cup \mathcal{T}_{M_1}(\Omega_{2,n,1})$, see Figure 7.

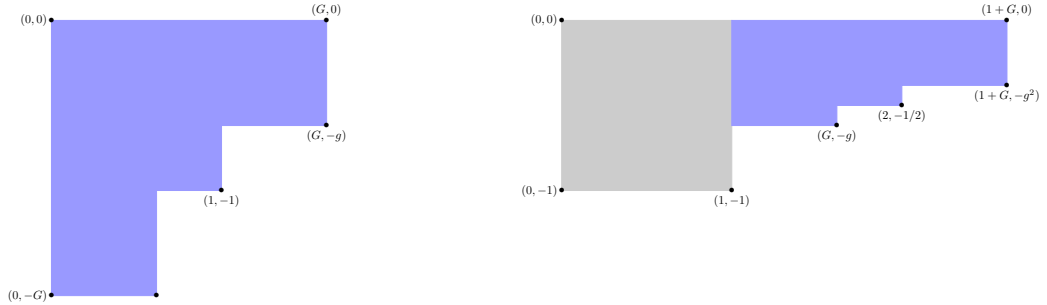


FIGURE 7. For each $m \geq 3$, the map $\mathcal{T}_{M_1}$, with $M_1 = \begin{pmatrix} 1 & 1 \\ 0 & 1 \end{pmatrix}$, sends the planar domain $\Omega_{2,n,1}$ to $\Omega_{3,n,1} \setminus \{(x, y) \mid x \geq 1\}$, see § 4.2. Here, $n=5$. On the left: The planar domain $\Omega_{2,5,1}$. On the right: $\Omega_{3,5,1}$ with the $\mathcal{T}_{M_1}$ -image of $\Omega_{2,5,1}$ highlighted.

4.3. Change of variable calculation. Due to the invariance of μ under $\mathcal{T}_M$ in general and thus under $\mathcal{T}_{M_1}$ in particular, for each $n \geq 3$ we have

$$\int_{\Omega_{3,n,1} \setminus ([0,1] \times [-1,0])} \tau_{3,n,1}(x) d\mu = \int_{\Omega_{2,n,1}} \tau_{3,n,1}(M_1 \cdot x) d\mu.$$

Now, for $x \geq 1$ one has $\tau_{3,n,1}(x) = \tau(A_3^k C_3^2, x)$ for some $k \geq 1$. Since both $\tau(A_3, x)$ and $\tau(A_2, x)$ vanish for any x , the cocycle property of $\tau(M, x)$ gives $\tau_{3,n,1}(M_1 \cdot x) = \tau(A_3^k C_3^2 M_1, x) = \tau(C_3^2 M_1, x) = \tau(C_2, x) = \tau(A_2^\ell C_2, x)$ for any $\ell \in \mathbb{Z}$ and any $0 \leq x < t_{2,n}$. Therefore,

$$\int_{\Omega_{3,n,1} \setminus ([0,1] \times [-1,0])} \tau_{3,n,1}(x) d\mu = \int_{\Omega_{2,n,1}} \tau_{2,n,1}(x) d\mu.$$

Finally, Theorem 5 gives

$$(32) \quad \int_{\Omega_{3,n,1} \setminus ([0,1] \times [-1,0])} \tau_{3,n,1}(x) d\mu = \text{vol}(T^1(G_{2,n} \setminus \mathbb{H})).$$

5. PROOF OF THE CONJECTURE

5.1. Conjecture reduction given in [CKS3]. By [[CKS3], Theorem 4], for each $n \geq 3$, the function $\alpha \mapsto h(T_{3,n,\alpha})\mu(\Omega_{3,n,\alpha})$ is constant on $(0, 1)$. Furthermore, for all such α , by [[CKS3], Lemma 149],

$$h(T_{3,n,\alpha})\mu(\Omega_{3,n,\alpha}) = \int_{\Omega_{3,n,\alpha}} \tau_{3,n,\alpha}(x) d\mu = \int_{\Omega_{3,n,0}} \tau_{3,n,0}(x) d\mu = \int_{\Omega_{3,n,1}} \tau_{3,n,1}(x) d\mu.$$

Thus, to prove the first statement of Theorem 1, it suffices to show that $\int_{\Omega_{3,n,1}} \tau_{3,n,1}(x) d\mu$ equals the volume of the unit tangent bundle of $G_{3,n} \setminus \mathbb{H}$.

As mentioned above, [[CKS2], Proposition 12] shows that $\Omega_{3,n,1}$ has the form of (31) and in particular contains the square $[0, 1] \times [-1, 0]$. By [[CKS3], Lemma 151],

$$(33) \quad \int_{[0,1] \times [-1,0]} -2 \log x d\mu = \pi^2/3.$$

As pointed out in the proof of [[CKS2], Proposition 12], for $x < 1$ one has $T_{3,n,1}(x) = A_3^k C_3 \cdot x$, with $k \in \mathbb{Z}$. Hence, for such x , $\tau_{3,n,1}(x) = -2 \log x$. It follows that

$$\int_{\Omega_{3,n,1}} \tau_{3,n,1}(x) d\mu = \pi^2/3 + \int_{\Omega_{3,n,1} \setminus ([0,1] \times [-1,0])} \tau_{3,n,1}(x) d\mu.$$

5.2. Final argument. From this last, (32) gives

$$(34) \quad \int_{\Omega_{3,n,1}} \tau_{3,n,1}(x) d\mu = \pi^2/3 + \text{vol}(T^1(G_{2,n} \setminus \mathbb{H})).$$

From (13), trivial algebra shows that $\text{vol}(T^1(G_{3,n} \setminus \mathbb{H})) = \pi^2/3 + \text{vol}(T^1(G_{2,n} \setminus \mathbb{H}))$. Combined with (34), this shows

$$\int_{\Omega_{3,n,1}} \tau_{3,n,1}(x) d\mu = \text{vol}(T^1(G_{3,n} \setminus \mathbb{H}))$$

and the conjecture holds!

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