# Fuchsian groups 

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## 1 Hyperbolic Geometry

### 1.1 Hyperbolic Metric (Lecture 1)

The hyperbolic plane is the metric space $(\mathbb{H}, \rho)$, where

$$
\mathbb{H}:=\{z \in \mathbb{C} \mid \operatorname{Im}(z)>0\}
$$

is the upper open halfplane of complex numbers and $\rho: \mathbb{H} \times \mathbb{H} \rightarrow \mathbb{R}$ is a certain metric, which we define below.

By definition, a $C^{1}$-curve in $\mathbb{H}$ is any $C^{1}$-map $z:[0,1] \rightarrow \mathbb{H}$. We write $z$ in the form

$$
z(t)=x(t)+i y(t),
$$

where $x$ and $y$ are real functions. Recall that the Euclidean length of $z$ is the number

$$
\ell(z):=\int_{0}^{1} \sqrt{\left(x^{\prime}(t)\right)^{2}+\left(y^{\prime}(t)\right)^{2}} d t=\int_{0}^{1}\left|z^{\prime}(t)\right| d t
$$

The hyperbolic length of $z$ is defined to be the number

$$
\begin{equation*}
h(z):=\int_{0}^{1} \frac{\sqrt{\left(x^{\prime}(t)\right)^{2}+\left(y^{\prime}(t)\right)^{2}}}{y(t)} d t=\int_{0}^{1} \frac{\left|z^{\prime}(t)\right|}{\operatorname{Im}(z(t))} d t . \tag{1.1.1}
\end{equation*}
$$

The hyperbolic distance from a point $z_{1} \in \mathbb{H}$ to a point $z_{2} \in \mathbb{H}$ is defined to be the number

$$
\begin{equation*}
\rho\left(z_{1}, z_{2}\right):=\inf h(z), \tag{1.1.2}
\end{equation*}
$$

where the infimum is taken over over all $C^{1}$-curves $z$ in $\mathbb{H}$ with $z(0)=z_{1}$ and $z(1)=z_{2}$.
Theorem 1.1.1 The function $\rho: \mathbb{H} \times \mathbb{H} \rightarrow \mathbb{R}$ is a metric on $\mathbb{H}$, i.e. for any $z_{1}, z_{2}, z_{3} \in \mathbb{H}$ we have
(1) $\rho\left(z_{1}, z_{1}\right)=0$ and $\rho\left(z_{1}, z_{2}\right)>0$ if $z_{1} \neq z_{2}$,
(2) $\rho\left(z_{1}, z_{2}\right)=\rho\left(z_{2}, z_{1}\right)$,
(3) $\rho\left(z_{1}, z_{2}\right) \leqslant \rho\left(z_{1}, z_{2}\right)+\rho\left(z_{2}, z_{3}\right)$.

Proof. (1) is clear, (3) will be proved in later (see Corollary 1.2.9). Now we prove (2). Let $z(t)=x(t)+i y(t), t \in[0,1]$ be a $C^{1}$-curve from $z_{1}$ to $z_{2}$. We define $s(t):=1-t$. Then

$$
w(t):=z(s(t))=x(s(t))+y(s(t)),
$$

$t \in[0,1]$ is $C^{1}$-curve from $z_{2}$ to $z_{1}$. In the following computation we use the chain rule for the differentiation and the substution rule for the integration.

$$
\begin{gathered}
h(w)=\int_{0}^{1} \frac{\mid\left(w^{\prime}(t) \mid\right.}{w(t)} d t=\int_{0}^{1} \frac{\left|z^{\prime}(s(t)) \cdot s^{\prime}(t)\right|}{z(s(t))} d t=-\int_{0}^{1} \frac{\left|z^{\prime}(s(t))\right|}{z(s(t))} \cdot s^{\prime}(t) d t \\
=-\int_{1}^{0} \frac{\left|z^{\prime}(s)\right|}{z(s)} d s=\int_{0}^{1} \frac{\left|z^{\prime}(s)\right|}{z(s)} d s=h(z) .
\end{gathered}
$$

### 1.2 The group of Möbius transformations

Let $A=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in \mathrm{SL}_{2}(\mathbb{C})$, i.e. $a, b, c, d \in \mathbb{C}, a d-b c=1$. The map

$$
\begin{aligned}
T_{A}: \mathbb{C} \cup\{\infty\} & \rightarrow \mathbb{C} \cup\{\infty\} \\
z & \mapsto \frac{a z+b}{c z+d}
\end{aligned}
$$

is called the Möbius transformation of $\mathbb{C} \cup\{\infty\}$ associated with $A$. Here we agree that

$$
\frac{a \infty+b}{c \infty+d}:=\left\{\begin{array}{l}
a / c, \text { of } c \neq 0 \\
\infty, \text { lf } c=0
\end{array}\right.
$$

Note that

$$
\begin{equation*}
T_{A} \circ T_{B}=T_{A B}, \quad T_{E}=i d, \tag{1.2.1}
\end{equation*}
$$

where $E$ is the identity matrix. The group

$$
\text { Möb }_{\mathbb{C}}:=\left\{T_{A} \mid A \in \mathrm{SL}_{2}(\mathbb{C})\right\}
$$

with respect to the composition is called the group of Möbius transformations of $\mathbb{C} \cup\{\infty\}$. The set

$$
\operatorname{Möb}_{\mathbb{R}}:=\left\{T_{A} \mid A \in \mathrm{SL}_{2}(\mathbb{R})\right\}
$$

forms a subgroup of $\mathrm{Möb}_{\mathbb{C}}$.
Theorem 1.2.1 Any transformation $T_{A} \in \mathrm{Möb}_{\mathbb{R}}$ maps the set $\mathbb{H}$ homeomorphically to itself.

Proof. We denote $w=T_{A}(z)$. Then we have

$$
\begin{gather*}
w=\frac{a z+b}{c z+d}=\frac{(a z+b)(c \bar{z}+d)}{|c z+d|^{2}}=\frac{a c|z|^{2}+a d z+b c \bar{z}+b d}{|c z+d|^{2}}, \\
\operatorname{Im}(w)=\frac{w-\bar{w}}{2 i}=\frac{z-\bar{z}}{2 i|c z+d|^{2}}=\frac{\operatorname{Im}(z)}{|c z+d|^{2}} . \tag{1.2.2}
\end{gather*}
$$

Hence, if $\operatorname{Im}(z)>0$ then $\operatorname{Im}(w)>0$. Therefore $T_{A}$ maps $\mathbb{H}$ to itself. The bijectivity of $\left(T_{A}\right)_{\mid \mathbb{H}}: \mathbb{H} \rightarrow \mathbb{H}$ follows from the formula $\left(T_{A}\right)_{\mid \mathbb{H}} \circ\left(T_{A^{-1}}\right)_{\mid \mathbb{H}}=i d_{\mid \mathbb{H}}$. The continuity of $\left(T_{A}\right)_{\mid \mathbb{H}}$ and of the inverse map $\left(T_{A^{-1}}\right)_{\mid \mathbb{H}}$ is obvious.

## Definition 1.2.2 .

1) $\left(\operatorname{Möb}_{\mathbb{R}}\right)_{\mid \mathbb{H}}$ denotes the set of all maps from Möb $_{\mathbb{R}}$ restricted to $\mathbb{H}$.
2) By $Z(G)$, we denote the center of the group $G$. We have

$$
Z\left(\mathrm{GL}_{2}(\mathbb{R})\right)=\left\{\left.\left(\begin{array}{cc}
r & 0 \\
0 & r
\end{array}\right) \right\rvert\, r \in \mathbb{R} \backslash\{0\}\right\}, \quad Z\left(\mathrm{SL}_{2}(\mathbb{R})\right)=\left\{ \pm\left(\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right)\right\}
$$

3) The groups

$$
\operatorname{PGL}_{2}(\mathbb{R}):=\mathrm{GL}_{2}(\mathbb{R}) / Z\left(\mathrm{GL}_{2}(\mathbb{R})\right) \quad \text { und } \quad \operatorname{PSL}_{2}(\mathbb{R}):=\mathrm{SL}_{2}(\mathbb{R}) / Z\left(\mathrm{SL}_{2}(\mathbb{R})\right)
$$

are called projective general linear groups and projective special linear groups.
Theorem 1.2.3 $\left(\mathrm{Möb}_{\mathbb{R}}\right)_{\mid \mathbb{H}} \cong \mathrm{PSL}_{2}(\mathbb{R})$.
Beweis. The map

$$
\begin{aligned}
\phi: \mathrm{SL}_{2}(\mathbb{R}) & \rightarrow\left(\mathrm{Möb}_{\mathbb{R}}\right)_{\mid \mathbb{H}} \\
A & \mapsto\left(T_{A}\right)_{\mid \mathbb{H}}
\end{aligned}
$$

is an epimorphism with $\operatorname{Ker} \phi=\left\{ \pm\left(\begin{array}{ll}1 & 0 \\ 0 & 1\end{array}\right)\right\}$. Therefore we have

$$
\left(\mathrm{Möb}_{\mathbb{R}}\right)_{\mid \mathbb{H}} \cong \mathrm{SL}_{2}(\mathbb{R}) / \operatorname{Ker} \phi \cong \mathrm{PSL}_{2}(\mathbb{R})
$$

Theorem 1.2.4 The transformations from $\left(\mathrm{Möb}_{\mathbb{R}}\right)_{\mathbb{H}}$ preserve the hyperbolic lengths of $C^{1}$-curves in $\mathbb{H}$.

Proof. Let

$$
\begin{aligned}
T: \mathbb{H} & \rightarrow \mathbb{H}, \\
z & \mapsto \frac{a z+b}{c z+d}
\end{aligned}
$$

be a transformation from $\left(\mathrm{Möb}_{\mathbb{R}}\right)_{\mid \mathbb{H}}$ with

$$
\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \in \mathrm{SL}_{2}(\mathbb{R})
$$

Let $z(t)=x(t)+i y(t)$ be a $C^{1}$-curve in $\mathbb{H}$. Denote $w(t):=T(z(t))$. We shall prove that $h(w)=h(z)$ :

$$
\begin{equation*}
h(w)=\int_{0}^{1} \frac{\left|w^{\prime}(t)\right|}{\operatorname{Im}(w(t))} d t=\int_{0}^{1} \frac{\left|T^{\prime}(z(t)) \cdot z^{\prime}(t)\right|}{\operatorname{Im}(T(z(t)))} d t \tag{1.2.3}
\end{equation*}
$$

We have

$$
T^{\prime}(z)=\frac{a(c z+d)-c(a z+b)}{(c z+d)^{2}}=\frac{1}{(c z+d)^{2}} .
$$

Using formula (1.2.2), we deduce

$$
\begin{equation*}
\left|T^{\prime}(z)\right|=\frac{\operatorname{Im}(T(z))}{\operatorname{Im}(z)} \tag{1.2.4}
\end{equation*}
$$

Substituting into (1.2.3), we obtain

$$
h(w)=\int_{0}^{1} \frac{\left|z^{\prime}(t)\right|}{\operatorname{Im}(z(t))} d t=h(z) .
$$

Corollary 1.2.5 The function $\rho$ is $\mathrm{Möb}_{\mathbb{R}}$-invariant, i.e. for any two points $z_{1}, z_{2} \in \mathbb{H}$ and any $T \in \operatorname{Möb}_{\mathbb{R}}$, we have

$$
\rho\left(z_{1}, z_{2}\right)=\rho\left(T\left(z_{1}\right), T\left(z_{2}\right)\right) .
$$

Proof. Using the definition of $\rho$ and Theorem 1.2.4, we have

$$
\rho\left(z_{1}, z_{2}\right)=\inf _{z \in \mathcal{Z}} h(z)=\inf _{z \in \mathcal{Z}} h(T(z))=\inf _{w \in \mathcal{W}} h(w)=\rho\left(T\left(z_{1}\right), T\left(z_{2}\right)\right) .
$$

Here $\mathcal{Z}$ is the set of all $C^{1}$-curves from $z_{1}$ to $z_{2}$, and $\mathcal{W}$ is the set of all $C^{1}$-curves from $T\left(z_{1}\right)$ to $T\left(z_{2}\right)$.

Definition 1.2.6 (a) For any $r \in \mathbb{R}$, let $\mathbf{A}_{r}$ be the open Euclidean axis in $\mathbb{H}$, whose closure begins at the point $r$ of the real axis and is perpendicular to this axis.
b) For any two different numbers $r_{1}, r_{2} \in \mathbb{R}$, let $\mathbf{C}_{r_{1}, r_{2}}$ be the open Euclidean half-circle in $\mathbb{H}$, whose closure begins at the point $r_{1}$ of the real axis and ends at the point $r_{2}$. These open axes and open half-circles are called geodesic lines in $\mathbb{H}$.
c) For any two different points $z_{1}, z_{2}$ in $\mathbb{H}$, we define the curve $\left[z_{1}, z_{2}\right]$ in $\mathbb{H}$ with the beginning $z_{1}$ and the end $z_{2}$ as follows:

Case 1. $\operatorname{Re}\left(z_{1}\right)=\operatorname{Re}\left(z_{2}\right)$.
Then there exists a unique axis $\mathbf{A}_{r}$ containing $z_{1}, z_{2}$. Let $\left[z_{1}, z_{2}\right]$ be the curve which goes "inside" of $\mathbf{A}_{r}$ from $z_{1}$ to $z_{2}$. More precisely, let

$$
\left[z_{1}, z_{2}\right](t)=(1-t) z+t z .
$$

Case 2. $\operatorname{Re}\left(z_{1}\right) \neq \operatorname{Re}\left(z_{2}\right)$.
Then there exists a unique half-circle $\mathbf{C}_{r_{1}, r_{2}}$, containing $z_{1}, z_{2}$. Let $\left[z_{1}, z_{2}\right]$ be the "natural" curve which goes "inside" of $\mathbf{C}_{r_{1}, r_{2}}$ from $z_{1}$ to $z_{2}$.
In both cases we call the curve $\left[z_{1}, z_{2}\right]$ the geodesic segment from $z_{1}$ to $z_{2}$.

## Lemma 1.2.7

(a) For any real $r$, there exists a transformation from $\mathrm{Möb}_{\mathbb{R}}$, which sends $\mathbf{A}_{r}$ to $\mathbf{A}_{0}$.
(b) For any two different reals $r_{1}, r_{2}$, there exists a transformation from Möb $\mathbb{R}_{\mathbb{R}}$, which sends $\mathbf{C}_{r_{1}, r_{2}}$ to $\mathbf{A}_{0}$.

Proposition 1.2.8 Let $z$ be a $C^{1}$-curve in $\mathbb{H}$ from $z_{1}$ to $z_{2}$. Then

$$
\begin{equation*}
h(z) \geqslant\left|\ln \frac{\operatorname{Im}\left(z_{2}\right)}{\operatorname{Im}\left(z_{1}\right)}\right| . \tag{1.2.5}
\end{equation*}
$$

In particular,

$$
\begin{equation*}
\rho\left(z_{1}, z_{2}\right) \geqslant\left|\ln \frac{\operatorname{Im}\left(z_{2}\right)}{\operatorname{Im}\left(z_{1}\right)}\right| . \tag{1.2.6}
\end{equation*}
$$

Proof. Denote $z(t)=x(t)+i y(t)$. Then
$h(z)=\int_{0}^{1} \frac{\sqrt{\left(x^{\prime}(t)\right)^{2}+\left(y^{\prime}(t)\right)^{2}}}{y(t)} d t \geqslant \int_{0}^{1} \frac{\left|y^{\prime}(t)\right|}{y(t)} d t \geqslant\left|\int_{0}^{1} \frac{y^{\prime}(t)}{y(t)} d t\right|=\left|\ln (y(t))_{0}^{1}\right|=\left|\ln \frac{y(1)}{y(0)}\right|=\left|\ln \frac{\operatorname{Im}\left(z_{2}\right)}{\operatorname{Im}\left(z_{1}\right)}\right|$.
The second formula follows from the first one; use the definition of $\rho$ as infimum of $h(z)$ over all $z$ from $z_{1}$ to $z_{2}$.

Theorem 1.2.9 Let $z_{1}, z_{2}$ be two different points in $\mathbb{H}$ and let $z$ be an arbitrary $C^{1}$-curve in $\mathbb{H}$ from $z_{1}$ to $z_{2}$. Then $h(z) \geqslant h\left(\left[z_{1}, z_{2}\right]\right)>0$.

Proof. By Lemma 1.2.7 and Theorem 1.2.4, we can assume that $z_{1}=i a$ and $z_{2}=i b$ $(b>a>0)$. Denote $z(t)=x(t)+i y(t)$. Then, using formula (1.2.5), we deduce

$$
h(z) \geqslant\left|\ln \frac{\operatorname{Im}\left(z_{2}\right)}{\operatorname{Im}\left(z_{1}\right)}\right|=\ln \frac{b}{a} .
$$

It suffices to prove that

$$
h\left(\left[z_{1}, z_{2}\right]\right)=\ln \frac{b}{a} .
$$

Denote $w=\left[z_{1}, z_{2}\right]$ and recall that $w(t)=i(a+t(b-a)), t \in[0,1]$. Then

$$
h(w)=\int_{0}^{1} \frac{\left|w^{\prime}(t)\right|}{\operatorname{Im}(w(t))} d t=\int_{0}^{1} \frac{b-a}{a+t(b-a)} d t=\left.\ln (a+t(b-a))\right|_{0} ^{1}=\ln b-\ln a=\ln \frac{b}{a} .
$$

## Corollary 1.2 .10

(1) For any two points $z_{1}, z_{2}$ in $\mathbb{H}$ holds

$$
\begin{equation*}
\rho\left(z_{1}, z_{2}\right)=h\left(\left[z_{1}, z_{2}\right]\right) . \tag{1.2.7}
\end{equation*}
$$

(2) For any two real numbers $b>a>0$, we have

$$
\begin{equation*}
\rho(i a, i b)=\ln \frac{b}{a} \tag{1.2.8}
\end{equation*}
$$

(3) For any three points $z_{1}, z_{2}, z_{3}$ in $\mathbb{H}$, we have

$$
\rho\left(z_{1}, z_{3}\right) \leqslant \rho\left(z_{1}, z_{2}\right)+\rho\left(z_{2}, z_{3}\right)
$$

The equality happens exactly in the case where $z_{2} \in\left[z_{1}, z_{3}\right]$.
Proof. Statement (1) follows from Theorem 1.2.9, statement (2) was established in the proof of this theorem. We prove (3). W.l.o.g., $z_{1}=i a, z_{3}=i b \quad(b>a>0)$. Then $\rho\left(z_{1}, z_{2}\right)+\rho\left(z_{2}, z_{3}\right) \geqslant\left|\ln \frac{\operatorname{Im}\left(z_{2}\right)}{\operatorname{Im}\left(z_{1}\right)}\right|+\left|\ln \frac{\operatorname{Im}\left(z_{3}\right)}{\operatorname{Im}\left(z_{2}\right)}\right| \geqslant \ln \frac{\operatorname{Im}\left(z_{2}\right)}{\operatorname{Im}\left(z_{1}\right)}+\ln \frac{\operatorname{Im}\left(z_{3}\right)}{\operatorname{Im}\left(z_{2}\right)}=\ln \frac{\operatorname{Im}\left(z_{3}\right)}{\operatorname{Im}\left(z_{1}\right)}=\ln \frac{b}{a}=\rho\left(z_{1}, z_{3}\right)$.

### 1.3 Some formulas for the hyperbolic metric $\rho$ (Lecture 2)

Lemma 1.3.1 We have $\mathrm{SL}_{2}(\mathbb{R})=\left\langle\left\{A_{r}, B_{r} \mid r \in \mathbb{R}\right\}\right\rangle=\left\langle\left\{A_{r} \mid r \in \mathbb{R}\right\} \cup\{C\}\right\rangle$, where

$$
A_{r}=\left(\begin{array}{cc}
1 & r \\
0 & 1
\end{array}\right), \quad B_{r}=\left(\begin{array}{cc}
1 & 0 \\
r & 1
\end{array}\right), \quad C=\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right)
$$

Corollary 1.3.2 We have $\operatorname{Möb}_{\mathbb{R}}=\left\langle\left\{\varphi_{r} \mid r \in \mathbb{R}\right\} \cup\{\psi\}\right\rangle$, where

$$
\begin{aligned}
\varphi_{r}: & z \rightarrow z+r, \\
\psi: & z \mapsto-\frac{1}{z} .
\end{aligned}
$$

Definition 1.3.3 The set $\widehat{\mathbb{C}}:=\mathbb{C} \cup\{\infty\}$ is called Riemannian sphere. The cross-ratio of four different points $z_{1}, z_{2}, z_{3}, z_{4} \in \widehat{\mathbb{C}}$ is defined to be

$$
\left(z_{1}, z_{2} ; z_{3}, z_{4}\right)=\frac{\left(z_{1}-z_{2}\right)\left(z_{3}-z_{4}\right)}{\left(z_{2}-z_{3}\right)\left(z_{4}-z_{1}\right)}
$$

If some $z_{i}$ equals to $\infty$, we first cancel the (two) terms containing $\infty$, using the rule

$$
\frac{\infty}{ \pm \infty}= \pm 1
$$

We also use some natural rules like

$$
\frac{a}{\infty}=0(a \in \mathbb{R}), \quad \frac{b}{0}=\infty(b \in \mathbb{R} \backslash\{0\})
$$

Theorem 1.3.4 The cross-ratio is $\mathrm{Möb}_{\mathbb{R}}$-invariant.
Proof. The proof follows with the help of Corollary 1.3.2.
Theorem 1.3.5 Let $z, w$ be two different points in $\mathbb{H}$. Let $z^{*}$ and $w^{*}$ be the ends of the geodesic line passing through $z$ and $w$. We assume that the order of the four points in the completion of this geodesic line is $z^{*}, z, w, w^{*}$ (see Fig.1). Then

$$
\rho(z, w)=\ln \left(w, z^{*} ; z, w^{*}\right) .
$$



Fig. 1.

Proof. By Lemma 1.2.7, there exist $T \in$ Möb $_{\mathbb{R}}$ which maps the geodesic line to the imaginary axis $\mathbf{A}_{0}$. Applying the maps $\theta_{k}$ and $\psi$, and using Theorem 1.3.4, we may assume that $T\left(z^{*}\right)=0, T\left(w^{*}\right)=\infty$. Then $T(z)=i a$ and $T(w)=i b$ for some $0<a<b$. By Corollary 1.2.10, we have

$$
\rho(z, w)=\rho(T(z), T(w))=\rho(i a, i b)=\ln \frac{b}{a} .
$$

We also have

$$
\left(w, z^{*} ; z, w^{*}\right)=\left(T(w), T\left(z^{*}\right) ; T(z), T\left(w^{*}\right)\right)=(i b, 0 ; i a, \infty)=\frac{(i b-0)(i a-\infty)}{(0-i a)(\infty-i b)}=\frac{b}{a}
$$

Definition 1.3.6 We define the following three functions from $\mathbb{R}$ to $\mathbb{R}$ :

$$
\begin{aligned}
\operatorname{ch}(t):=\frac{e^{t}+e^{-t}}{2} & \text { (hyperbolic cosinus), } \\
\operatorname{sh}(t):=\frac{e^{t}-e^{-t}}{2} & \text { (hyperbolic sinus), } \\
\operatorname{th}(t):=\frac{\operatorname{sh}(t)}{\operatorname{ch}(t)}=\frac{e^{t}-e^{-t}}{e^{t}+e^{-t}} & \text { (hyperbolic tangens) }
\end{aligned}
$$

Theorem 1.3.7 For every two points $z, w \in \mathbb{H}$ the following formulas are valid.

$$
\begin{aligned}
& \text { (1) } \begin{array}{l}
\rho(z, w) \\
=\ln \frac{|z-\bar{w}|+|z-w|}{|z-\bar{w}|-|z-w|} ; \\
\text { (2) } \operatorname{ch} \rho(z, w) \\
=1+\frac{|z-w|^{2}}{2 \operatorname{Im}(z) \operatorname{Im}(w)} ; \\
\text { (3) } \operatorname{sh}\left[\frac{1}{2} \rho(z, w)\right]
\end{array}=\frac{|z-w|}{2(\operatorname{Im}(z) \operatorname{Im}(w))^{1 / 2}} ; \\
& \text { (4) } \operatorname{ch}\left[\frac{1}{2} \rho(z, w)\right]
\end{aligned}=\frac{|z-\bar{w}|}{2(\operatorname{Im}(z) \operatorname{Im}(w))^{1 / 2}} ;
$$

Proof. One can directly prove that these equations are equivalent. Therefore, we prove only (3).

By Theorem 1.2.4, the left side of (3) is $\mathrm{Möb}_{\mathbb{R}}$-invariant. With the help of Corollary 1.3.2, we first verify that the right side of (3) is also $\mathrm{Möb}_{\mathbb{R}^{-}}$-invariant. Since the right side of (3) is evidently $\varphi_{r}$-invariant, it suffices to check that it is $\psi$-invariant:

$$
\frac{|\psi(z)-\psi(w)|}{2(\operatorname{Im}(\psi(z)) \operatorname{Im}(\psi(w)))^{1 / 2}}=\frac{\left|\frac{-1}{z}-\frac{-1}{w}\right|}{2\left(\operatorname{Im}\left(\frac{-1}{z}\right) \operatorname{Im}\left(\frac{-1}{w}\right)\right)^{1 / 2}}=\frac{\left|\frac{-1}{z}-\frac{-1}{w}\right|}{2\left(\frac{\operatorname{Im}(z)}{|z|^{2}} \frac{\operatorname{Im}(w)}{|w|^{2}}\right)^{1 / 2}}=\frac{|z-w|}{2(\operatorname{Im}(z) \operatorname{Im}(w))^{1 / 2}} .
$$

Therefore, after application of a suitable $T \in \operatorname{Möb}_{\mathbb{R}}$, we may assume that $z=i a, w=i b$ $(a<b)$. Using $\rho(i a, i b)=\ln \frac{b}{a}$, we can easily verify (3).

### 1.4 Isometries of $\mathbb{H}$

Definition 1.4.1 A map $f: \mathbb{H} \rightarrow \mathbb{H}$ is called an isometry if $\rho\left(f\left(z_{1}\right), f\left(z_{2}\right)\right)=\rho\left(z_{1}, z_{2}\right)$ for all $z_{1}, z_{2} \in \mathbb{H}$.

The set of all Isometries the hyperbolic plane $\mathbb{H}$ is a group. ${ }^{1}$ This group is denoted by $\operatorname{Isom}(\mathbb{H})$.

Lemma 1.4.2 The following statements hold.
(1) $\left(\text { Möb }_{\mathbb{R}}\right)_{\mid \mathbb{H}} \leqslant \operatorname{Isom}(\mathbb{H})$.
(2) Isometries map geodesic lines to geodesic lines (see Definition 1.2.6).

Proof. (1) follows from Theorem 1.2.4. The proof of (2) can be extracted from the proof of Theorem 1.2.10.

Notation 1.4.3 We will use the following maps from $\mathrm{Möb}_{\mathbb{R}}$ :

$$
\theta_{k}: z \mapsto k z \quad\left(k \in \mathbb{R}_{+}\right) .
$$

We introduce the map from $\mathbb{H}$ to $\mathbb{H}$ which is not from $\mathrm{Möb}_{\mathbb{R}}$ (show this!).

$$
\eta: z \mapsto-\bar{z} .
$$

Theorem 1.4.4 We have

$$
\begin{aligned}
\operatorname{Isom}(\mathbb{H}) & =\left\langle\left(\operatorname{Möb}_{\mathbb{R}}\right)_{\mid \mathbb{H}}, \eta\right\rangle \\
& \cong \operatorname{PSL}_{2}(\mathbb{R}) \rtimes \mathbb{Z}_{2} .
\end{aligned}
$$

Beweis. We prove $\operatorname{Isom}(\mathbb{H}) \leqslant\left\langle\left(\operatorname{Möb}_{\mathbb{R}}\right)_{\mathbb{H}}, \eta\right\rangle$. Let $\varphi \in \operatorname{Isom}(\mathbb{H})$. Then $\varphi$ maps geodesic lines to geodesic lines. Let $I:=\{i r \mid r>0\}$. Then $\varphi(I)$ is a geodesic line. By Lemma 1.2.7, there exists $T \in\left(\mathrm{Möb}_{\mathbb{R}}\right)_{\mid \mathbb{H}}$ such that $T \circ \varphi(I)=I$. Using the maps $z \mapsto k z(k>0)$ and $z \mapsto-\frac{1}{z}$ if necessary, we may assume that $T \circ \varphi$ fixes the point $i$ and maps the axes $(0, i]$ and $[i, \infty)$ onto itself. From this, it follows that the isometry $T \circ \varphi$ fixes all points on $I$. We prove that $T \circ \varphi \in\{i d, \eta\}$. Let $z=x+i y \in \mathbb{H}$ and let $T \circ \varphi(z)=u+i v$. Then for all $t>0$ :

$$
\rho(x+i y, i t)=\rho(T \circ \varphi(z), T \circ \varphi(i t))=\rho(u+i v, i t) .
$$

By Theorem 1.3.7 (3), we have

$$
\frac{x^{2}+(y-t)^{2}}{t y}=\frac{u^{2}+(v-t)^{2}}{t v}
$$

It follows that $v=y$ and $x= \pm u$. Thus, $T \circ \varphi(z) \in\{z,-\bar{z}\}$. Since every isometry is continuous, the map $T \circ \varphi$ is either the identity on $\mathbb{H}$, or $\eta$.

[^0]
### 1.5 Hyperbolic Area (Lecture 3)

Definition 1.5.1 Let $A$ be an open set in $\mathbb{H}$. The hyperbolic area of $A$ is defined by the formula

$$
\mu(A)=\int_{A} \frac{d x d y}{y^{2}}
$$

In this integral, we consider $A$ as a subset of $\mathbb{R}^{2}$.
Remark 1.5.2 (Transformation-formula)

- Let $A$ be an open subset of $\mathbb{R}^{2}$.
- Let $\varphi: A \rightarrow \mathbb{R}^{2}$ be an injective differentiable function with continuous partial derivatives. We write $\varphi=\left(\varphi_{1}, \varphi_{2}\right)$ mit $\varphi_{i}: A \rightarrow \mathbb{R}$ and use the Jacobian of $\varphi$ :

$$
J(\varphi)=\operatorname{det}\left(\begin{array}{ll}
\frac{\partial \varphi_{1}}{\partial x} & \frac{\partial \varphi_{1}}{\partial y} \\
\frac{\partial \varphi_{2}}{\partial x} & \frac{\partial \varphi_{2}}{\partial y}
\end{array}\right) .
$$

- Let $f: \varphi(A) \rightarrow \mathbb{R}$ be a continuous function. Then the function $f$ has an integral on $\varphi(A)$ if and only if the function $(f \circ \varphi) \cdot|J(\varphi)|$ has an integral on $A$. In this case, we have

$$
\int_{\varphi(A)} f(x, y) d x d y=\int_{A}(f \circ \varphi)(x, y) \cdot|J(\varphi)(x, y)| d x d y
$$

Theorem 1.5.3 The function $\mu$ is $\operatorname{Möb}(\mathbb{R})$-invariant, i.e. for every open subset $A \subseteq \mathbb{H}$ and every $T \in\left(\mathrm{Möb}_{\mathbb{R}}\right)_{\mid \mathbb{H}}$ we have

$$
\mu(T(A))=\mu(A) .
$$

Proof. Let

$$
T(z)=\frac{a z+b}{c z+d} \quad(a, b, c, d \in \mathbb{R}, a d-b c=1)
$$

We write $z=x+i y$. Then there exist functions $u, v$ with $T(z)=u(x, y)+i v(x, y)$. Since $T$ is complex-differentiable (holomorph), $u$ and $v$ satisfy Cauchy-Riemann equations

$$
\begin{aligned}
& \frac{\partial u}{\partial x}=\frac{\partial v}{\partial y}, \\
& \frac{\partial u}{\partial y}=-\frac{\partial v}{\partial x} .
\end{aligned}
$$

Using these equations, we compute the Jacobian of the map

$$
\begin{gathered}
\varphi:(x, y) \mapsto(u(x, y), v(x, y)) . \\
J(\varphi):=\operatorname{det}\left(\begin{array}{cc}
\frac{\partial u}{\partial x} & \frac{\partial u}{\partial y} \\
\frac{\partial v}{\partial x} & \frac{\partial v}{\partial y}
\end{array}\right)=\left(\frac{\partial u}{\partial x}\right)^{2}+\left(\frac{\partial v}{\partial x}\right)^{2}=\left|\left(\frac{\partial u}{\partial x}+i \frac{\partial v}{\partial x}\right)\right|^{2}=\left|\frac{d T}{d z}\right|^{2}=\frac{1}{|c z+d|^{4}} .
\end{gathered}
$$

Then ${ }^{2}$

$$
\begin{gathered}
\mu(T(A))=\int_{T(A)} \frac{d x d y}{(\operatorname{Im}(z))^{2}}=\int_{A} \frac{1}{(\operatorname{Im}(T(z)))^{2}} \cdot|J(T)| d x d y \\
\stackrel{(1.2 .2)}{=} \int_{A} \frac{|c z+d|^{4}}{(\operatorname{Im}(z))^{2}} \cdot \frac{1}{|c z+d|^{4}} d x d y=\mu(A) .
\end{gathered}
$$

### 1.6 Angles in $\mathbb{H}$

Definition 1.6.1 Let $\gamma_{1}:\left[c_{1}, d_{1}\right] \rightarrow \mathbb{H}$ and $\gamma_{2}:\left[c_{2}, d_{2}\right] \rightarrow \mathbb{H}$ be two injective diffrentiable curves in $\mathbb{H}$, which pass through some common point $z$. The hyperbolic angle between $\gamma_{1}$ and $\gamma_{2}$ at the point $z$ is the Euclidian angle between the tangent lines $\zeta_{1}$ and $\zeta_{2}$ to these curves at the point $z$ :

$$
\angle\left(\gamma_{1}, \gamma_{2} ; z\right):=\angle_{\mathrm{e}}\left(\zeta_{1}, \zeta_{2} ; z\right) .
$$

Theorem 1.6.2 The transformations $T \in\left(\operatorname{Möb}_{\mathbb{R}}\right)_{\mid \mathbb{H}}$ are conform, i.e. they preserve the orientation and angles between $C^{1}$-curves:

$$
\angle\left(T\left(\gamma_{1}\right), T\left(\gamma_{2}\right) ; T(z)\right)=\angle\left(\gamma_{1}, \gamma_{2} ; z\right)
$$

### 1.7 Gauß-Bonnet formula

A hyperbolic n-gon in $\mathbb{H}$ is a closed subset of $\mathbb{H}$, which is bounded by $n$ hyperbolic segments of the form $[z, w]$. We also consider hyperbolic n-gons in the extension $\overline{\mathbb{H}}:=\mathbb{H} \cup \mathbb{R} \cup\{\infty\}$.

Theorem 1.7.1 (Gauß-Bonnet) Let $\Delta$ be a hyperbolic triangle in $\overline{\mathbb{H}}$ with the angles $\alpha, \beta, \gamma$. Then

$$
\mu(\Delta)=\pi-\alpha-\beta-\gamma
$$

Proof. Let $\Delta=A B C$.
Case 1. Let $A, B \in \mathbb{H}$ and $C \in \mathbb{R} \cup\{\infty\}$.
Since Möbius transformations preserve areas and angles, we may assume that $C=\infty$.
Then the side $A B$ lies on the halfcircle $\mathbf{C}_{r_{1}, r_{2}}$. We may assume that 0 is the center of $\mathbf{C}_{r_{1}, r_{2}}$. Let $R$ be the radius of this halfcircle. Sine the sides $A C$ and $B C$ are geodesic lines and $C=\infty$, they are vertical axes. Let $a$ and $b$ be $x$-coordinates of these axes. Then

$$
\mu(\Delta)=\int_{\Delta} \frac{d x d y}{y^{2}}=\int_{a}^{b} d x \int_{\sqrt{R^{2}-x^{2}}}^{\infty} \frac{d y}{y^{2}}=\int_{a}^{b} \frac{d x}{\sqrt{R^{2}-x^{2}}}
$$

After substitution $x=R \cos \theta$, we obtain

$$
\mu(\Delta)=\int_{\pi-\alpha}^{\beta} \frac{-\sin \theta d \theta}{\sin \theta}=\pi-\alpha-\beta
$$

[^1]Case 2. Let $A, B, C \in \mathbb{H}$.
Case 3. Let $A \in \mathbb{H}$ und $B, C \in \mathbb{R} \cup\{\infty\}$.
Case 4. Let $A, B, C \in \mathbb{R} \cup\{\infty\}$.
These cases can be reduced to Case 1 .

### 1.8 Hyperbolic trigonometry

Theorem 1.8.1 Let $\Delta$ be a geodesic triangle in $\mathbb{H}$ with finite hyperbolic lengths $a, b, c$ of its sides and with the non-zero angles $\alpha, \beta, \gamma$ (opposite to the corresponding sides). Then
(1) $\frac{\operatorname{sh} a}{\sin \alpha}=\frac{\operatorname{sh} b}{\sin \beta}=\frac{\operatorname{sh} c}{\sin \gamma} . \quad$ (sinus theorem)
(2) $\operatorname{ch} c=\operatorname{ch} a \operatorname{ch} b-\operatorname{sh} a \operatorname{sh} b \cos \gamma . \quad$ (first cosinus theorem)
(3) $\operatorname{ch} c=\frac{\cos \alpha \cos \beta+\cos \gamma}{\sin \alpha \sin \beta}$. (second cosinus theorem)

Theorem 1.8.2 If two geodesic triangles in $\mathbb{H}$ have the same angles, then there exists an isometry which maps one triangle to the other.

Theorem 1.8.3 (Pythagoras Theorem for $\mathbb{H})$ Let $\Delta$ be a geodesic triangle in $\mathbb{H}$ with finite hyperbolic lengths $a, b, c$ of its sides and with the non-zero angles $\alpha, \beta, \gamma$ (opposite to the corresponding sides). If $\gamma=\frac{\pi}{2}$, then

$$
\operatorname{ch} c=\operatorname{ch} a \operatorname{ch} b
$$

## 2 Fuchsian groups

### 2.1 Classification of elements of $\mathrm{PSL}_{2}(\mathbb{R})$

Definition 2.1.1 Let $A=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right)$ be an element of the group $\mathrm{SL}_{2}(\mathbb{R})$ different from $\pm E$. The number

$$
\operatorname{Tr}(A):=a+d
$$

is called the trace of $A$.

- $A$ is called elliptic if $|\operatorname{Tr}(A)|<2$.
- $A$ is called parabolic if $|\operatorname{Tr}(A)|=2$.
- $A$ is called hyperbolic if $|\operatorname{Tr}(A)|>2$.

A nontrivial element from $\mathrm{PSL}_{2}(\mathbb{R})$ is called elliptic, parabolic, or hyperbolic if some (equivalently any) its preimage in $\mathrm{SL}_{2}(\mathbb{R})$ is elliptic, parabolic, or hyperbolic, respectively.

We write $A \underset{\mathrm{SL}_{2}(\mathbb{R})}{\sim} B$ if $A$ and $B$ are conjugate by a matrix from $\mathrm{SL}_{2}(\mathbb{R})$.
Lemma 2.1.2 Suppose that $\pm E \neq A=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in \mathrm{SL}_{2}(\mathbb{R})$, and let $\lambda_{1}, \lambda_{2}$ be eigenvectors of $A$. Then the following holds.

1) $A$ is hyperbolic if and only if $\lambda_{1}, \lambda_{2} \in \mathbb{R}, \lambda_{1} \neq \lambda_{2}$. In this case

$$
A \underset{\operatorname{SL}_{2}(\mathbb{R})}{\sim}\left(\begin{array}{cc}
\lambda & 0 \\
0 & 1 / \lambda
\end{array}\right)
$$

for some $\lambda \in \mathbb{R}, \lambda \neq \pm 1$.
2) $A$ is parabolic if and only if $\lambda_{1}=\lambda_{2} \in\{-1,1\}$. In this case

$$
A_{\mathrm{SL}_{2}(\mathbb{R})}^{\sim} \pm\left(\begin{array}{ll}
1 & \alpha \\
0 & 1
\end{array}\right)
$$

for some $\alpha \in\{-1,1\}$.
3) $A$ is elliptic if and only if $\lambda_{1}, \lambda_{2} \in \mathbb{C} \backslash \mathbb{R}$. In this case $\lambda_{2}=\bar{\lambda}_{1},\left|\lambda_{1}\right|=\left|\lambda_{2}\right|=1$ and

$$
A \underset{\mathrm{SL}_{2}(\mathbb{R})}{\sim}\left(\begin{array}{cc}
\cos \theta & \sin \theta \\
-\sin \theta & \cos \theta
\end{array}\right)
$$

for some $\theta \in(0, \pi) \cup(\pi, 2 \pi)$.
Definition 2.1.3 The element $A=\left[\begin{array}{ll}a & b \\ c & d\end{array}\right] \in \operatorname{PSL}_{2}(\mathbb{R})$ acts on $\mathbb{H}$ as the map $z \mapsto \frac{a z+b}{c z+d}$. The set of fixed points of $A$ in $\mathbb{H}$ is

$$
\operatorname{Fix}(A):=\left\{z \in \mathbb{H} \left\lvert\, z=\frac{a z+b}{c z+d}\right.\right\}
$$

The action of $A$ on $\mathbb{H}$ can be naturally extended to an action of $A$ on the compactified hyperbolic plane $\widehat{\mathbb{H}}:=\mathbb{H} \cup \partial \mathbb{H}$, where $\partial \mathbb{H}=\mathbb{R} \cup\{\infty\}$. The subset $\partial \mathbb{H}$ is called the boundary of $\mathbb{H}$. The set of fixed points of $A$ in $\widehat{\mathbb{H}}$ is denoted by $\widehat{\operatorname{Fix}}(A)$.

Theorem 2.1.4 For any $A \in \mathrm{PSL}_{2}(\mathbb{R})$ the following holds.

1) If $A$ is hyperbolic, then $\widehat{\operatorname{Fix}}(A)$ consists of two points in $\partial \mathbb{H}$. One of them is attracting and the other one is repelling.
2) If $A$ is parabolic, then $\widehat{\operatorname{Fix}}(A)$ consists of a single point in $\partial \mathbb{H}$.
3) If $A$ elliptic, then $\widehat{\operatorname{Fix}}(A)$ consists of a single point in $\mathbb{H}$.

Definition 2.1.5 Suppose that $A \in \mathrm{PSL}_{2}(\mathbb{R})$ is hyperbolic. The geodesic line in $\mathbb{H}$ connecting two fixed points of $A$ is called the axis of $A$ and is denoted by $\operatorname{Axis}(A)$.

Bemerkung 2.1.6 The axis of a hyperbolic element $A \in \mathrm{PSL}_{2}(\mathbb{R})$ is $A$-invariant.
Bemerkung 2.1.7 If we know the type of $A$ (hyperbolic, parabolic, or elliptic) and the fixed points of $A$, we can describe the action of $A$ on $\widehat{\mathbb{H}}$ on a qualitative level, i.e. with the help of pictures.

Bemerkung 2.1.8 The element $A=\left[\begin{array}{ll}a & b \\ c & d\end{array}\right] \in \mathrm{PSL}_{2}(\mathbb{R})$ fixes $\infty$ if and only if $c=0$. In this case the corresponding transformation $\varphi_{A}: \mathbb{C} \rightarrow \mathbb{C}$ has the form $z \mapsto a^{2} z+b a$ and we have the following.

- If $a= \pm 1$, then $A$ is parabolic.
- If $a \neq \pm 1$, then $A$ is hyperbolic with $\operatorname{Fix}(A)=\left\{\infty, \frac{b a}{1-a^{2}}\right\}$.


## Subsections 2.2 and 2.3 contain necessary information about topo-

 logical spaces and topological groups.
### 2.2 Topological spaces

Definition 2.2.1 (Topology, topological space, open and closed sets)
Let $X$ be a set. A topology $\mathfrak{T}$ on $X$ is a set of some subsets of $X$ (each such subset is called an open set in $X$ ), which satisfies the following axioms:
(1) The empty set $\varnothing$ and the set $X$ are open.
(2) The intersection of finitely many open sets is open.
(3) the union of arbitrary set of open sets is open

The pair $(X, \mathfrak{T})$ is called a topological space. Sometimes we simply write $X$ for the topological space if the topology $\mathfrak{T}$ is defined. A subset $U$ of $X$ is called closed if $X \backslash U$ is open.

Definition 2.2.2 (Basis of a topology)
Let $(X, \mathfrak{T})$ be a topological space. A subset $B \subseteq \mathfrak{T}$ is called a basis of the topology $\mathfrak{T}$ if each open set of $X$ is a union of some open sets belonging to $B$.

Definition 2.2.3 (Neighborhood) Let $(X, \mathfrak{T})$ be a topological space and let $x$ be a point of $X$. A subset $U$ of $X$ is called a neighborhood of $x$ if there exists an open subset $\mathcal{O}$ such that $x \in \mathcal{O} \subseteq U$.

## Definition 2.2.4 (Continuous maps)

Let $\left(X_{1}, \mathfrak{T}_{1}\right)$ and $\left(X_{2}, \mathfrak{T}_{2}\right)$ be two topological spaces. A map $f: X_{1} \rightarrow X_{2}$ is called continuous if for every open set $U$ in $X_{2}$ the full preimage $f^{-1}(U)$ is open in $X_{1}$.

Definition 2.2.5 (Induced topologie und quotient topology)
(1) Let $(X, \mathfrak{T})$ be a topological space and let $Y$ be a subset of $X$. We define on $Y$ the induced topology $\mathfrak{T}_{Y}$ as follows: A subset $S \subseteq Y$ is defined to be open in $Y$ if there exists an open set $\mathcal{O}$ in $X$ such that $S=\mathcal{O} \cap Y$.
(2) Let $(X, \mathfrak{T})$ be a topological space and let $Y$ be a set. Let $f: X \rightarrow Y$ be a map. We define on $Y$ the quotient topology as follows: A subset $U \subseteq Y$ is defined to be open in $Y$ if the full preimage $f^{-1}(U)$ is open in $X$.

## Remark 2.2.6

(1) Let $(X, \mathfrak{T})$ be a topological space and let $Y$ be a subset of $X$. The inclusion map $i: Y \rightarrow X, y \mapsto y$, becomes continuous if we endow $Y$ with the induced topology. Moreover, the induced topology on $Y$ is the weakest topology on $Y$ for which the map $i$ is continuous.
(2) Let $(X, \mathfrak{T})$ be a topological space and let $Y$ be a set. Let $f: X \rightarrow Y$ be a map. The map $f: X \rightarrow Y$ becomes continuous if we endow $Y$ with the quotient topology. Moreover, the quotient topology on $Y$ is the weakest topology on $Y$ for which the map $f$ is continuous.

Definition 2.2.7 (Discrete topological space and a discrete subset of a topological space)
(1) A topological space $(X, \mathfrak{T})$ is called discrete if one from two equivalent statements is valid:
(a) For every point $x \in X$ the set $\{x\}$ is open.
(b) Every subset of $X$ is open.
(2) A subset $X$ of a topological space $Y$ is called discrete if $X$ with the induced topology is discrete. Equivalently, if for every point $x \in X$ there exists an open set $\mathcal{O}_{x}$ in $Y$ such that $\mathcal{O}_{x} \cap X=\{x\}$.

Definition 2.2.8 (Hausdorff space) A topological space ( $X, \mathfrak{T}$ ) is called Hausdorff space if for every two different points $x_{1}$ and $x_{2}$ in $X$ there exist open sets $\mathcal{O}_{1}$ and $\mathcal{O}_{2}$ such that

$$
x_{1} \in \mathcal{O}_{1}, x_{2} \in \mathcal{O}_{2} \text { and } \mathcal{O}_{1} \cap \mathcal{O}_{2}=\varnothing
$$

Definition 2.2.9 (First countable topological spaces) A topological space ( $X, \mathfrak{T}$ ) is called first countable if the following holds:

For every point $x \in X$, there exists a countable collection of neihborhoods $U_{1}, U_{2}, \ldots$ of $x$ such that, for every neighborhood $U$ of $x$, there exists $i \in \mathbb{N}$ such that $U_{i} \subseteq U$.

Remark 2.2.10 (Continuity and countable sequences)
(1) Let $(X, \mathfrak{T})$ be a topological space and let $\left(Y, \mathfrak{T}_{Y}\right)$ be a subspace with the induced topology. Then the following holds:
(i) If $(X, \mathfrak{T})$ is first countable, then $\left(Y, \mathfrak{T}_{Y}\right)$ is first countable.
(ii) Let $\left(y_{n}\right)_{n \in \mathbb{N}} \subseteq Y$ be a sequence and $y \in Y$ a point. Then $y_{n} \rightarrow y$ in $\left(Y, \mathfrak{T}_{Y}\right)$ if and only if $y_{n} \rightarrow y$ in $(X, \mathfrak{T})$.
(2) Let $\left(X_{1}, \mathfrak{T}_{1}\right)$ and $\left(X_{2}, \mathfrak{T}_{2}\right)$ be two topological spaces. Let $f: X_{1} \rightarrow X_{2}$ be a map. We will compare the following two conditions.
(a) The map $f$ is continuous.
(b) For every sequence $\left(x_{n}\right)_{n \in \mathbb{N}}$ in $X_{1}$ and every $x \in X_{1}$ the following holds: If $x_{n} \rightarrow x$ in $\left(X_{1}, \mathfrak{T}_{1}\right)$, then $f\left(x_{n}\right) \rightarrow f(x)$ in $\left(X_{2}, \mathfrak{T}_{2}\right)$.

Condition (b) follows from condition (a), but not conversely, in general. If $\left(X_{i}, \mathfrak{T}_{i}\right), i=1,2$, are first countable then (a) and (b) are equivalent.

Remark 2.2.11 (Criterium of discreteness) Let $(X, \mathfrak{T})$ be a first countable topological space. Then this space is discrete if for every point $x \in X$ and every sequence $\left(x_{n}\right)_{n \in \mathbb{N}} \subseteq X$ with $x_{n} \rightarrow x$ we have $x_{n}=x$ for all sufficiently large $n$.

Definition 2.2.12 (Compact spaces and compact subsets)
(1) A topological space $(X, \mathfrak{T})$ is called compact if for each covering $X=\underset{i \in I}{\cup} U_{i}$ with $U_{i} \in \mathfrak{T}$, there exists a finite subset $I_{0} \subseteq I$ such that $X=\cup_{i \in I_{0}} U_{i}$.
(2) A subset $Y$ of a topological space $(X, \mathfrak{T})$ is called compact if one of the following equivalent conditions is satisfied:
(a) $\left(Y, \mathfrak{T}_{Y}\right)$ is a compact space.
(b) For each covering $Y \subseteq \bigcup_{i \in I} U_{i}$ with $U_{i} \in \mathfrak{T}$, there exists a finite subset $I_{0} \subseteq I$ such that $Y \subseteq \cup_{i \in I_{0}} U_{i}$.

Remark 2.2.13 (Compactness for metric spaces)
Let $(X, d)$ be a metric space. For $r>0$ and $x \in X$ the set

$$
B_{r}(x):=\left\{x^{\prime} \in X \mid d\left(x, x^{\prime}\right)<r\right\}
$$

is called an open ball with the center $x$ and radius $r$. A subset $Y$ of $X$ is called bounded if $Y$ lies in some open ball.

Let $\mathfrak{T}$ be the topology on $X$ for which the set $\left\{B_{r}(x) \mid x \in X, r>0\right\}$ is a basis. Then
(1) The topological space $(X, \mathfrak{T})$ is Hausdorff and first countable.
(2) A subset $Y$ in $\mathbb{R}^{n}$ is compact if and only if $Y$ is closed and bounded.

Lemma 2.2.14 Every discrete and closed subset $M$ of a compact topological space $X$ is finite.

Proof. For every $m \in M$ let $O(m)$ be an open neighborhood of $m$ in $X$ satisfying $O(m) \cap M=\{m\}$. Then

$$
X=(X \backslash M) \cup(\underset{m \in M}{\cup} O(m))
$$

is a covering of $X$ by open sets. Since $X$ is compact, there exists a finite subset $M_{0} \subseteq M$ such that

$$
X=(X \backslash M) \cup\left(\underset{m \in M_{0}}{\cup} O(m)\right)
$$

Then

$$
M=\left(\underset{m \in M_{0}}{\cup} O(m)\right) \cap M=\underset{m \in M_{0}}{\cup}(O(m) \cap M)=M_{0} .
$$

Definition 2.2.15 (Product topology)
Let $\left(X_{1}, \mathfrak{T}_{1}\right)$ and $\left(X_{2}, \mathfrak{T}_{2}\right)$ be two topological spaces. The product topology $\mathfrak{T}$ on $X_{1} \times X_{2}$ consists of all possible unions of sets of the form $U \times V$, where $U \in \mathfrak{T}_{1}$ and $V \in \mathfrak{T}_{2}$.

Remark 2.2.16 The product topology is the weakest topology on $X_{1} \times X_{2}$, for which the projections $\mathrm{pr}_{1}: X_{1} \times X_{2} \rightarrow X_{1}$ and $\mathrm{pr}_{2}: X_{1} \times X_{2} \rightarrow X_{2}$ are continuous.

### 2.3 Topological groups

Definition 2.3.1 A group $G$, which is simultaneously a topological space, is called a topological group, if the maps $\cdot: G \times G \rightarrow G,(x, y) \mapsto x y$ and ${ }^{-1}: G \rightarrow G, x \mapsto x^{-1}$ are continuous.

Remark 2.3.2 Let $G$ be a topological group. Then the following statements are valid.
(1) For every open set $U \subseteq G$ and any element $g \in G$ the sets $g U$ and $U g$ are open.
(2) For every neighborhood $V$ of 1 in $G$, there exists a neighborhood $U$ of 1 in $G$ such that $U U^{-1} \subseteq V$.

Proof. (2) Since $1 \cdot 1=1$ and since the multiplication in $G$ is continuous, there exist two neighborhoods of 1 , say $U_{1}, U_{2}$ such that $U_{1} U_{2} \subseteq V$. For $U_{3}=U_{1} \cap U_{2}$ we have $U_{3} U_{3} \subseteq V$. For $U=U_{3} \cap U_{3}^{-1}$ we finally have $U U^{-1} \subseteq V$.

Definition 2.3.3 A subgroup $H$ of a topological group $G$ is called discrete in $G$ if $H$ is discrete as a subset of the topological space $G$.

Theorem 2.3.4 Every discrete subgroup $H$ of a Hausdorff topological group $G$ is closed.
Proof. Since $H$ is discrete in $G$, there exists an open neighborhood $V$ of 1 such that $V \cap H=\{1\}$. Then there exists an open neighborhood $U$ of 1 such that $U U^{-1} \subseteq V$.

We show that $G \backslash H$ is open. Let $g \in G \backslash H$. We shall show that there exists an open neighborhood $W$ of $g$ which does not contain elements of $H$. Try $W=U g$. Suppose it contains some $h \in H$. Then, since $G$ is Hausdorff, we can find an open neighborhood $W_{1} \subseteq W$ of $g$, which does not contain $h$. Suppose that $W_{1}$ contains some other $h_{1} \in H$. Then $h g^{-1}, h_{1} g^{-1} \in U$. Then $h h_{1}^{-1}=\left(h g^{-1}\right)\left(h_{1} g^{-1}\right)^{-1} \in U U^{-1} \subseteq V$, hence $h h_{1}^{-1}=1$, a contradiction. Thus, $W_{1}$ is an open neighborhood of $g$ that does not contain elements of $H$. Hence, $G \backslash H$ is open.

Remark 2.3.5 The variant of Theorem 2.3.4 for subspaces of topological spaces is not valid. Indeed, consider the interval $[0,1]$ as a topological space with the topology that is induced by the canonical topology on $\mathbb{R}$. Then the topological space $[0,1]$ is Hausdorff and compact. Die subset $\left\{\left.\frac{1}{n} \right\rvert\, n \in \mathbb{N}\right\}$ of the topological space $[0,1]$ is discrete but not closed.

Corollary 2.3.6 Every discrete subgroup $H$ of a Hausdorff kompact topological group $G$ is finite.

Proof. The proof follows straightforwardly from Theorem 2.3.4 and Lemma 2.2.14.
Corollary 2.3.7 Every discrete subgroup of the orthogonal group $O(n)$ is finite.
Lemma 2.3.8 (Criterium of discreteness of a topological subgroup) Let $G$ be a topological group with a first countable Hausdorff topology. Then the following holds:

A subgroup $H$ of $G$ is discrete in $G$ if and only if from $h_{n} \rightarrow e$ (where $h_{n} \in H$ and $e$ is a neutral element) follows that $h_{n}=e$ for all sufficiently large $n$.

### 2.4 First two definitions of a Fuchsian group

The group $\mathrm{SL}_{2}(\mathbb{R})$ can be considered as a subset of $\mathbb{R}^{4}$ by identifying

$$
\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \rightsquigarrow(a, b, c, d)
$$

Thus, $\mathrm{SL}_{2}(\mathbb{R})$ can be considered as a topological group (and even a metric space) with the topology induced from $\mathbb{R}^{4}$. Now we consider the canonical epimorphism

$$
\psi: \mathrm{SL}_{2}(\mathbb{R}) \rightarrow \mathrm{PSL}_{2}(\mathbb{R})
$$

From now on, we consider $\mathrm{PSL}_{2}(\mathbb{R})$ as a topological group with the quotient topology determined by $\psi$.

Definition 2.4.1 (first definition of a Fuchsian group) A Fuchsian group is a discrete subgroup of $\mathrm{PSL}_{2}(\mathbb{R})$.

To understand the discreteness in $\operatorname{PSL}_{2}(\mathbb{R})$ better, we need the following two general statements.

Lemma 2.4.2 Let $G_{1}$ be a topological group, $G_{2}$ be a group and $\varphi: G_{1} \rightarrow G_{2}$ an epimorphism. We consider $G_{2}$ as a topological group with respect to the quotient topology. Let $H$ be a subgroup of $G_{2}$. Then the following statements are valid.
(a) If $\varphi^{-1}(H)$ is discrete in $G_{1}$, then $H$ is discrete in $G_{2}$.
(b) Suppose additionally that $G_{1}$ is Hausdorff and that $\operatorname{ker}(\varphi)$ is finite.

If $H$ is discrete in $G_{2}$, then $\varphi^{-1}(H)$ is discrete in $G_{1}$.
Lemma 2.4.3 Let $G_{1}$ be a topological group, $G_{2}$ be a group. Let $\varphi: G_{1} \rightarrow G_{2}$ be an epimorphism and $G_{2}$ is endowed by the quotient topology. Then the following statements are valid.
(1) If a subset $\mathcal{O} \subseteq G_{1}$ is open, then its image $\varphi(\mathcal{O})$ is open.
(2) If $G_{1}$ is first countable, then $G_{2}$ is also first countable.

Corollary 2.4.4 The topological group $\mathrm{PSL}_{2}(\mathbb{R})$ is Hausdorff and first countable.
Using this corollary and the discreteness criterium 2.3.8, we give the second (equivalent) definition of a Fuchsian group.

Definition 2.4.5 (second definition of a Fuchsian group) A subgroup $H \leqslant \mathrm{PSL}_{2}(\mathbb{R})$ is called Fuchsian if for any sequence $\left(h_{n}\right)_{n \in \mathbb{N}}$ of elements of $H$ with $h_{n} \rightarrow 1$, there exists $n_{0} \in \mathbb{N}$ such that for every $n \geqslant n_{0}$, we have $h_{n}=1$.

### 2.5 Proper discontinuous actions of groups on metric spaces

Definition 2.5.1 Let $X$ be a topological space and $S$ a subset of $X$. A point $x \in X$ is called an accumulation point of $S$ if every neighborhood of $x$ contains a point of $S$ different from $x$. The set of all accumulation points of $S$ in $X$ is denoted by $\mathbf{A P}_{X}(S)$.

Assumption. From now on we assume that $(X, d)$ is a metric space and $G$ is a group acting on $X$ by isometries, i.e. for any element $g \in G$ and any two points $x_{1}, x_{2} \in X$ we have

$$
d\left(x_{1}, x_{2}\right)=d\left(g\left(x_{1}\right), g\left(x_{2}\right)\right) .
$$

An action (by isometries) of $G$ on $X$ is denoted by $G \curvearrowright X$.
Definition 2.5.2 An action $G \curvearrowright X$ is called proper discontinuous (abbreviated PDA), if for every point $x \in X$ there exists a neighborhood $V$ of $x$ such that

$$
\begin{equation*}
|\{g \in G \mid g(V) \cap V \neq \emptyset\}|<\infty \tag{2.5.1}
\end{equation*}
$$

Lemma 2.5.3 A group $G$ acts on a metric space $X$ properly discintinuously if and only if the folowing two conditions are satisfied.
(1) For any $x \in X$ the orbit $G(x)$ does not have an accumulation point in $X$.
(2) For any $x \in X$ the stabilizer $\operatorname{St}_{G}(x):=\{g \in G \mid g(x)=x\}$ is finite.

$$
\text { Proof. }(1) \&(2) \Rightarrow(\mathrm{PDA}):
$$

Let $x \in X$ be an arbitrary point. It follows from (1) that there exists $\varepsilon>0$ with

$$
B_{\varepsilon}(x) \cap G(x)=\{x\} .
$$

We claim that $V:=B_{\varepsilon / 2}(x)$ satisfies (2.5.1). Indeed, if $g \in G$ is an element satisfying

$$
\begin{equation*}
g\left(B_{\varepsilon / 2}(x)\right) \cap B_{\varepsilon / 2}(x) \neq \emptyset \tag{2.5.2}
\end{equation*}
$$

then $d(x, g(x))<\varepsilon$, hence

$$
g(x) \in B_{\varepsilon}(x) \cap G(x)=\{x\}
$$

i.e. $g \in \operatorname{St}_{G}(x)$. $\mathrm{By}(2)$, this stabilizer is finite. Thus, there exists only finitely many $g \in G$ satisfying (2.5.2), and (PDA) is proved.
$\urcorner(2) \Rightarrow\urcorner$ (PDA) is evident. Now we prove
$\urcorner(1) \Rightarrow\urcorner(\mathrm{PDA}):$
By $\urcorner(1)$, there exists $x \in X$ such that $G(x)$ has an accumulation point $y$. Then there exist different $g_{i} \in G, i \in \mathbb{N}$, such that $g_{i}(x) \in B_{1 / i}(y) \backslash\{y\}$. In particular,

$$
g_{i}^{-1}\left(B_{1 / i}(y)\right) \cap g_{j}^{-1}\left(B_{1 / j}(y)\right) \neq \emptyset
$$

for all $i, j$. Therefore, for $j \geqslant i$ we have

$$
g_{j} g_{i}^{-1}\left(B_{1 / i}(y)\right) \cap B_{1 / i}(y) \neq \emptyset .
$$

Let $\varepsilon>0$. We take $i_{0} \in \mathbb{N}$ such that $1 / i_{0}<\varepsilon$. Then

$$
g_{j} g_{i_{0}}^{-1}\left(B_{\varepsilon}(y)\right) \cap B_{\varepsilon}(y) \neq \emptyset .
$$

for all $j \geqslant i_{0}$. This implies $\urcorner$ (PDA).
Lemma 2.5.4 Let $z_{0} \in \mathbb{H}$ and let $K$ be a compact subset of $\mathbb{H}$. Then the set

$$
M:=\left\{T \in \mathrm{SL}_{2}(\mathbb{R}) \mid T\left(z_{0}\right) \in K\right\}
$$

is compact.
Proof. It suffices to show that $M$ (as a subset of $\mathbb{R}^{4}$ ) is closed and bounded.

1) We prove that $M$ is closed. Consider the map

$$
\begin{aligned}
\psi: \mathrm{SL}_{2}(\mathbb{R}) & \rightarrow \mathbb{H}, \\
T & \mapsto T\left(z_{0}\right) .
\end{aligned}
$$

Then

$$
M=\left\{T \in \mathrm{SL}_{2}(\mathbb{R}) \mid \psi(T) \in K\right\}=\psi^{-1}(K)
$$

Since $\psi$ is continuous and $K$ is closed, $M$ is closed as well.
2) We prove that $M$ is bounded.
a) Since $K$ is bounded, there exists a constat $C_{1}>0$ such that $|z|<C_{1}$ for all $z \in K$. Then we have

$$
\begin{equation*}
\left|\frac{a z_{0}+b}{c z_{0}+d}\right|<C_{1} \tag{2.5.3}
\end{equation*}
$$

for all $T=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right)$ from $M$.
b) Since $K$ is bounded, there exists a constant $C_{2}>0$ such that

$$
\operatorname{Im}\left(\frac{a z_{0}+b}{c z_{0}+d}\right) \geqslant C_{2}
$$

By (1.2.2), we have

$$
\frac{\operatorname{Im}\left(z_{0}\right)}{\left|c z_{0}+d\right|^{2}} \geqslant C_{2}
$$

Therefore

$$
\begin{equation*}
\left|c z_{0}+d\right| \leqslant \sqrt{\frac{\operatorname{Im}\left(z_{0}\right)}{C_{2}}} \tag{2.5.4}
\end{equation*}
$$

This and (2.5.3) imply

$$
\begin{equation*}
\left|a z_{0}+b\right| \leqslant C_{1} \sqrt{\frac{\operatorname{Im}\left(z_{0}\right)}{C_{2}}} \tag{2.5.5}
\end{equation*}
$$

From (2.5.4) and (2.5.5) we deduce that $|a|,|b|,|c|,|d|$ are bounded.
Theorem 2.5.5 A subgroup $G$ of $\operatorname{PSL}_{2}(\mathbb{R})$ is Fuchsian if and only if $G$ acts properly discontinuously on $\mathbb{H}$.

Proof. Let $\psi: \mathrm{SL}_{2}(\mathbb{R}) \rightarrow \mathrm{PSL}_{2}(\mathbb{R})$ be the canonical epimorphism. Then (tutorial)

- $G$ is discrete in $\mathrm{PSL}_{2}(\mathbb{R})$ if and only if $\psi^{-1}(G)$ discrete in $\mathrm{SL}_{2}(\mathbb{R})$.
- $G$ acts on $\mathbb{H}$ totally discontinuously if and only if $\psi^{-1}(G)$ acts on $\mathbb{H}$ totally discontinuously.

Therefore it suffices to prove the following:
$A$ subgroup $G$ of $\mathrm{SL}_{2}(\mathbb{R})$ is Fuchsian if and only if $G$ acts properly discontinuously on $\mathbb{H}$.

1) Suppose that $G$ is discrete in $\mathrm{SL}_{2}(\mathbb{R})$.

To the contrary, suppose that the action $G \curvearrowright \mathbb{H}$ is not proper discontinuous. Then there exists $z_{0} \in \mathbb{H}$ such that the set

$$
G_{0}:=\left\{g \in G \mid g\left(B_{1}\left(z_{0}\right)\right) \cap B_{1}\left(z_{0}\right) \neq \emptyset\right\}
$$

is infinite. We set $K:=\overline{B_{2}\left(z_{0}\right)}$. Then $g\left(z_{0}\right) \in K$ for all $g \in G_{0}$. Therefore

$$
G_{0} \subseteq\left\{g \in \mathrm{SL}_{2}(\mathbb{R}) \mid g\left(z_{0}\right) \in K\right\} \cap G
$$

- The set $\left\{g \in \mathrm{SL}_{2}(\mathbb{R}) \mid g\left(z_{0}\right) \in K\right\}$ is compact (see Lemma 2.5.4).
- Since $G$ is discrete in $\mathrm{SL}_{2}(\mathbb{R}), G$ is closed (see Lemma 2.3.4).
- The intersection of a compact set and a discrete closed set is finite (see Lemma 2.2.14).

Then $G_{0}$ is finite. A contradiction.
2) Suppose $G$ is not discrete in $\mathrm{SL}_{2}(\mathbb{R})$. Then there exists an infinite sequence $\left(g_{k}\right)_{k \in \mathbb{N}}$ of different and nontrivial elements of $G$ such thatz $g_{k} \rightarrow e$. Each element from $\mathrm{SL}_{2}(\mathbb{R})$ fixes at most one point of $\mathbb{H}$. Let $z_{0} \in \mathbb{H}$ be a point, which is not fixed by any $g_{k}$. Then

- $g_{k}\left(z_{0}\right) \neq z_{0}$ for all $k$,
- $g_{k}\left(z_{0}\right) \rightarrow z_{0}$.

Thus, $z_{0}$ is an accumulation point of $G\left(z_{0}\right)$. By Lemma 2.5.3 the action $G \curvearrowright \mathbb{H}$ is not proper discontinuous.

Corollary 2.5.6 (Criterium for Fuchsian groups) A subgroup $G \leqslant \mathrm{PSL}_{2}(\mathbb{R})$ is Fuchsian if and only if for every point $z \in \mathbb{H}$ the orbit $G(z)$ does not have accumulation points in $\mathbb{H}$.

Proof. $(\Rightarrow)$ follows from Theorem 2.5.5 and Lemma 2.5.3.
$(\Leftarrow)$ Suppose that $G$ is not discrete. As in the second part of the proof of Theorem 2.5.5, there exists a point $z_{0} \in \mathbb{H}$ such that the orbit $G\left(z_{0}\right)$ has an accumulation point in $\mathbb{H}$. A contradiction.

Definition 2.5.7 The limit set of a subgroup $G \leqslant \operatorname{PSL}_{2}(\mathbb{R})$ is the set

$$
\Lambda(G)=\bigcup_{z \in \mathbb{H}} \mathbf{A P}_{\widehat{\mathbb{H}}}(G(z)) .
$$

Thus, this is the set of all accumulation points (in $\widehat{\mathbb{H}}$ ) of all orbits $G(z), z \in \mathbb{H}$.
Lemma 2.5.8 If $G$ is a Fuchsian group, then
(1) $\Lambda(G) \subseteq \mathbb{R} \cup\{\infty\}$,
(2) $G(\Lambda(G))=\Lambda(G)$.

Proof. (1) follows from Corollary 2.5.6, (2) from Definition 2.5.7.

## Examples.

1) For $G=\langle A\rangle$ with $A=\left[\begin{array}{cc}2 & 0 \\ 0 & 1 / 2\end{array}\right]$ we have $\Lambda(G)=\{0, \infty\}$.
2) For $G=\mathrm{PSL}_{2}(\mathbb{Z})$ we have $\Lambda(G)=\mathbb{R} \cup\{\infty\}$.

Corollary 2.5.9 For every Fuchsian group $G$ the set of fixpoints of all its elliptic elements does not have accumulation points in $\mathbb{H}$. With other words, the set

$$
\{z \in \mathbb{H} \mid g(z)=z \text { for some nontrivial } g \in G\}
$$

does not have accumulation points in $\mathbb{H}$.
Proof. The proof follows from Theorem 2.5.5 and Definition 2.5.2.

### 2.6 Some algebraic properties of Fuchsian groups

Recall our agreements:

$$
\frac{a \cdot 0+b}{c \cdot 0+d}:=\left\{\begin{array}{l}
b / d, \text { if } d \neq 0, \\
\infty, \text { if } d=0
\end{array} \quad \frac{a \infty+b}{c \infty+d}:=\left\{\begin{array}{l}
a / c, \text { if } c \neq 0, \\
\infty, \text { if } c=0
\end{array}\right.\right.
$$

The centralizer of an element $a$ of $G$ is the subgroup

$$
C_{G}(a):=\{g \in G \mid g a=a g\} .
$$

The normalizer of a subgroup $H$ of $G$ is the subgroup

$$
N_{G}(H):=\{g \in G \mid g H=H g\} .
$$

Lemma 2.6.1 Let $T, S$ be two nontrivial elements of Möb $_{\mathbb{R}}$. If $T S=S T$, then

$$
S(\widehat{\operatorname{Fix}}(T))=\widehat{\operatorname{Fix}}(T)
$$

Proof. Let $z \in \widehat{\operatorname{Fix}}(T)$. Then $S(z)=S T(z)=T S(z)$, hence $S(z) \in \widehat{\operatorname{Fix}}(T)$. Thus, $S(\widehat{\operatorname{Fix}}(T)) \subseteq \widehat{\operatorname{Fix}}(T)$. Similarly, $T S^{-1}=S^{-1} T$ implies $S^{-1}(\widehat{\operatorname{Fix}}(T)) \subseteq \widehat{\operatorname{Fix}}(T)$.

Theorem 2.6.2 Let $T, S$ be two nontrivial elements of $\mathrm{Möb}_{\mathbb{R}}$. Then

$$
T S=S T \Leftrightarrow \widehat{\operatorname{Fix}}(T)=\widehat{\operatorname{Fix}}(S)
$$

In particular, the types of $S$ and $T$ are coincide.
Proof. Consider three cases.

1) Let $T$ be parabolic. W.l.o.g. (and using conjugation), we may assume $T: z \mapsto z+1$. In particular, $\widehat{\operatorname{Fix}}(T)=\{\infty\}$.

- First suppose that $T S=S T$. Then, by Lemma 2.6.1, we have $S(\infty)=\infty$, hence $S: z \mapsto a z+b$. Furthermore, from $S T=T S$, we deduce that $a=1$, so $S: z \mapsto z+b$. Thus, $\widehat{\operatorname{Fix}}(S)=\{\infty\}$.
- Now suppose that $\widehat{\operatorname{Fix}}(T)=\widehat{\operatorname{Fix}}(S)$. Then $\widehat{\operatorname{Fix}}(S)=\{\infty\}$. Then $S: z \mapsto a z+b$, $a \neq 0$. If $a \neq 1$, then we have an additional fixed point. Thus, $S: z \mapsto z+b$ that implies $T S=S T$.

2) Let $T$ be hyperbolic. W.l.o.g. (and using conjugation), we may assume $T: z \mapsto k z$, $k>0, k \neq 1$. In particular, $\widehat{\operatorname{Fix}}(T)=\{0, \infty\}$.

- First suppose that $T S=S T$. Then, by Lemma 2.6.1, we have $S(\{0, \infty\})=\{0, \infty\}$. The case $S(0)=\infty$ and $S(\infty)=0$ is impossible, otherwise $S: z \mapsto-\mu / z$ for some $\mu>0$, and hence $T S \neq S T$. Thus, $S(0)=0$ and $S(\infty)=\infty$, i.e. $\widehat{\operatorname{Fix}}(S)=\widehat{\operatorname{Fix}}(T)$.
- Now suppose that $\widehat{\operatorname{Fix}}(S)=\widehat{\operatorname{Fix}}(T)$. Then $\widehat{\operatorname{Fix}}(S)=\{0, \infty\}$. Hence $S: z \mapsto \mu z$ for some $\mu>0$. Then $S T=T S$.

3) Let $T$ be elliptic. W.l.o.g. (and using conjugation), we may assume $T: z \mapsto$ $\frac{\cos \theta \cdot z+\sin \theta}{-\sin \theta \cdot z+\cos \theta}$ for some $\theta$. In particular, $\widehat{\operatorname{Fix}}(T)=\{i\}$.

- First suppose that $T S=S T$. Then, by Lemma 2.6.1, we have $i \in \widehat{\operatorname{Fix}}(S)$. Then $\{i\}=\widehat{\operatorname{Fix}}(S)$ by the classification Lemma 2.1.2.
- Now suppose that $\widehat{\operatorname{Fix}}(S)=\widehat{\operatorname{Fix}}(T)$. Then $\widehat{\operatorname{Fix}}(S)=\{i\}$. Hence, $z \mapsto \frac{\cos \varphi \cdot z+\sin \varphi}{-\sin \varphi \cdot z+\cos \varphi}$ for some $\varphi$. Then $S T=T S$.

Corollary 2.6.3 The centralizer in $\mathrm{Möb}_{\mathbb{R}}$ of a hyperbolic, parabolic, or elliptic element consists of id and all hyperbolic, parabolic, or elliptic elements, respectively, which have the same fixed points.

Corollary 2.6.4 Two hyperbolic elements commute if and only if they have the same axes.

Corollary 2.6.5 Let $A, B, C$ three nontrivial Möbius transformations. If $A B=B A$ and $B C=C B$, then $A C=C A$.

Lemma 2.6.6 Any discrete subgroup of $(\mathbb{R},+)$ is isomorphic to $\mathbb{Z}$. Any discrete subgroup of $S^{1}=\left([0,2 \pi],+{ }_{\bmod 2 \pi}\right)$ is isomorphic to $\mathbb{Z}_{n}$ for some finite $n$.

Theorem 2.6.7 Let $G$ be a Fuchsian group such that all nontrivial elements of $G$ have the same fixed points. Then $G$ is cyclic. Moreover, if $G$ contains a hyperbolic or a parabolic element, then $G \cong \mathbb{Z}$, and if $G$ contains an elliptic element, then $G \cong \mathbb{Z}_{n}$ for some $n$.

Proof. Let $g_{0} \in G \backslash\{1\}$ be a fixed element and $g \in G \backslash\{1\}$ an arbitrary. Consider three cases.

1) $g_{0}$ is hyperbolic. After conjugation, we may assume $g_{0}: z \mapsto k z$ for some $k>0$. Then $\widehat{\operatorname{Fix}}(g)=\operatorname{Fix}\left(g_{0}\right)=\{0, \infty\}$. Therefore $g: z \mapsto \mu z$ for some $\mu>0$, and we have

$$
G \leqslant\{z \mapsto \lambda z \mid \lambda>0\} \cong\left(\mathbb{R}_{+}, \cdot\right) \cong(\mathbb{R},+)
$$

By Lemma 2.6.6, $G \cong \mathbb{Z}$.
2) $g_{0}$ is parabolic. After conjugation, we may assume $g_{0}: z \mapsto z+1$ for some $k>0$. Then $\widehat{\operatorname{Fix}}(g)=\operatorname{Fix}\left(g_{0}\right)=\{\infty\}$. Therefore $g: z \mapsto z+b$ for some $b \in \mathbb{R}$, and we have

$$
G \leqslant\{z \mapsto z+\lambda \mid \lambda \in \mathbb{R}\} \cong(\mathbb{R},+) .
$$

By Lemma 2.6.6, $G \cong \mathbb{Z}$.
3) $g_{0}$ is elliptic. After conjugation, we may assume $g_{0}: z \mapsto \frac{\cos \theta \cdot z+\sin \theta}{-\sin \theta \cdot z+\cos \theta}$ for some $\theta$. Then $\widehat{\operatorname{Fix}}(g)=\operatorname{Fix}\left(g_{0}\right)=\{i\}$. Therefore

$$
G \leqslant\left\{\left.z \mapsto \frac{\cos \varphi \cdot z+\sin \varphi}{-\sin \varphi \cdot z+\cos \varphi} \right\rvert\, 0 \leqslant \varphi<2 \pi\right\} \cong S^{1}
$$

By Lemma 2.6.6, $G \cong \mathbb{Z}_{n}$.

Theorem 2.6.8 (1) Every abelian Fuchsian group is cyclic.
(2) If $G$ is Fuchsian and $g \in G \backslash\{1\}$, then the centralizer $C_{G}(g)$ is cyclic.

Proof. Statement (1) follows from Theorems 2.6.2 and 2.6.7. Statement (2) follows from Corollary 2.6.5 and statement (1).

Theorem 2.6.9 Let $G$ be a Fuchsian group. If $G$ is noncyclic, then $N=N_{\mathrm{PSL}_{2}(\mathbb{R})}(G)$ is also a Fuchsian group.

Proof. Suppose that $N$ is not discrete. Then there exists a sequence $\left(T_{i}\right)_{i \in \mathbb{N}}, T_{i} \rightarrow 1$, where all $T_{i}$ are different and nontrivial. Then for any $g \in G$ we have

$$
T_{i}^{-1} g T_{i} \rightarrow g
$$

Since $G$ is discrete, there exists $i_{0}=i_{0}(g)$ such that for any $i \geqslant i_{0}$ we have

$$
T_{i}^{-1} g T_{i}=g .
$$

Let $g_{1}, g_{2}$ be two nontrivial elements of $G$. Then there exists $T_{k}$ from the above sequence such that $T_{k}^{-1} g_{1} T_{k}=g_{1}$ and $T_{k}^{-1} g_{2} T_{k}=g_{2}$. By Corollary 2.6.5, $g_{1}$ and $g_{2}$ are commute, hence $G$ is abelian. By Theorem 2.6.8, $G$ is cyclic.

### 2.7 Elementary Fuchsian groups

Definition 2.7.1 A subgroup $G \leqslant \mathrm{PSL}_{2}(\mathbb{R})$ is called elementary if there exists a point $z \in \widehat{\mathbb{H}}$ such that the orbit $G(z)$ is finite.

Remark 2.7.2 Since $\mathbb{H}$ and $\mathbb{R} \cup\{\infty\}$ are $\operatorname{PSL}_{2}(\mathbb{R})$-invariant, we have

$$
G(z) \subseteq \mathbb{H} \text { or } G(z) \subseteq \mathbb{R} \cup\{\infty\}
$$

Theorem 2.7.3 Suppose that $G$ is a subgroup of $\operatorname{PSL}_{2}(\mathbb{R})$ such that all nontrivial elements of $G$ are elliptic. Then all elements of $G$ have a common fixed point in $\widehat{\mathbb{H}}$. In particular, $G$ is abelian and elementary.

Beweis. Let $A$ be an element in $G \backslash\{e\}$. After an appropriate conjugation, we have

$$
A=\left[\begin{array}{cc}
\cos \theta & -\sin \theta \\
\sin \theta & \cos \theta
\end{array}\right] .
$$

Let $B=\left[\begin{array}{ll}a & b \\ c & d\end{array}\right]$ be an arbitrary element of $G$. We have $\operatorname{det}(B)=1$. Then we have

$$
\begin{aligned}
\operatorname{Tr}\left(A B A^{-1} B^{-1}\right) & =2 a d \cos ^{2}(\theta)+\left(a^{2}+b^{2}+c^{2}+d^{2}\right) \sin ^{2}(\theta)-2 b c \cos ^{2}(\theta) \\
& =2 \cos ^{2}(\theta)+\left(a^{2}+b^{2}+c^{2}+d^{2}\right) \sin ^{2}(\theta) \\
& =2+\left(a^{2}+b^{2}+c^{2}+d^{2}-2\right) \sin ^{2}(\theta) \\
& =2+\left((a-d)^{2}+(b+c)^{2}\right) \sin ^{2}(\theta) \geqslant 2
\end{aligned}
$$

Therefore the commutator $[A, B]$ is either trivial or elliptic. But since it is not elliptic by assumption, we have $[A, B]=1$. Therefore $G$ is abelian. By Theorem 2.6.2, all nontrivial elements of $G$ have the same fixed points in $\widehat{\mathbb{H}}$. In particular, $G$ is elementary. E

Corollary 2.7.4 Every Fuchsian group $G$, whose nontrivial elements are elliptic is a finite cyclic group.

Proof. By Theorem 2.7.3, $G$ is abelian. By Theorem 2.6.2, the fixed points of nontrivial elements of $G$ coincide. By Theorem 2.6.8, $G$ is a finite cyclic group.

Theorem 2.7.5 Every elementary Fuchsian group $G$ is isomorphic to $\mathbb{Z}, \mathbb{Z}_{n}$ (for some finite $n$ ), or $D_{\infty}$ (the infinite dihedral group). Moreover, the following holds:
(1) If $G \cong \mathbb{Z}$, then $G$ is conjugate to the group $\langle z \mapsto k z\rangle$ for some $k>0$, or to the $\operatorname{group}\langle z \mapsto z+k\rangle$ for some $k>0$.
(2) If $G \cong \mathbb{Z}_{n}$, then $G$ is conjugate to the group

$$
\left\langle z \mapsto \frac{\cos \frac{2 \pi}{n} \cdot z+\sin \frac{2 \pi}{n}}{-\sin \frac{2 \pi}{n} \cdot z+\cos \frac{2 \pi}{n}}\right\rangle
$$

for some $n \in \mathbb{N}$.
(3) If $G \cong D_{\infty}$, then $G$ is conjugate to the group $H_{k}=\left\langle\theta_{k}, \psi\right\rangle$ for some $k>0$, where $\theta_{k}: z \mapsto k z$ and $\psi: z \mapsto-\frac{1}{z}$.

Proof. Let $\mathcal{O}$ be a finite orbit of $G$ in $\widehat{\mathbb{H}}$.
Case 1. Let $|\mathcal{O}|=1$.
Then $\mathcal{O}=\{\alpha\}$ for some $\alpha \in \widehat{\operatorname{Fix}}(G)$.
Fall 1.1. Sei $\alpha \in \mathbb{H}$.
Then all elements of $G \backslash\{e\}$ are elliptic. By Corollary 2.7.4, $G$ is a finite cyclic group.

Fall 1.2. Let $\alpha \in \mathbb{R} \cup\{\infty\}$.
Then $G$ has no elliptic elements. There are three subcases:
(a) $G$ contains both, hyperbolic and parabolic elements.

After an appropriate conjugation $G$ contains an element $g: z \mapsto \lambda z$, where $\lambda>0$. We have $\alpha \in \widehat{\operatorname{Fix}}(G) \subseteq \widehat{\operatorname{Fix}}(g)=\{0, \infty\}$. Therefore $\alpha=0$ or $\alpha=\infty$. If $\alpha=0$, we consider $\psi G \psi^{-1}$ instead of $G$ (recall that $\psi: z \rightarrow-\frac{1}{z}$ ). Then

- $\widehat{\operatorname{Fix}}\left(\psi G \psi^{-1}\right)=\psi(\widehat{\operatorname{Fix}}(G))=\psi(\alpha)=\{\infty\}$.
- $g=\psi g^{-1} \psi^{-1} \in \psi G \psi^{-1}$.

Therefore we may assume that $\alpha=\infty$. If necessary, we also can replace $g$ by $g^{-1}$ and assume that $\lambda>1$. Let $h$ be a parabolic element from $G$. From $\alpha \in \widehat{\operatorname{Fix}}(G) \subseteq \widehat{\operatorname{Fix}}(h)$ we deduce $\{\infty\}=\widehat{\operatorname{Fix}}(h)$. Then $h: z \rightarrow z+b$ for some $b \in \mathbb{R}$. We have

$$
g^{-n} h g^{n}(z)=z+\lambda^{-n} b
$$

Since $\lambda>0$, we have

$$
g^{-n} h g^{n} \rightarrow \text { id. }
$$

that contradicts the discreteness of $G$.
(b) $G \backslash\{e\}$ contains only parabolic elements.

Then each of them has $\alpha$ as a single fixed point. Then, by Theorem 2.6.7, we have $G \cong \mathbb{Z}$.
(c) $G \backslash\{e\}$ contains only hyperbolic elements.

As above, we may assume that $G$ contains $g: z \mapsto \lambda z(\lambda \neq 1)$. and $\alpha=\infty$. If $\langle g\rangle=G$, then we are done.
If $\langle g\rangle \neq G$, then we consider some $h \in G \backslash\langle g\rangle$. Since $\alpha \in \widehat{\operatorname{Fix}}(h)$, we have $h: z \mapsto a z+b$ for some $a \neq 0$ and $b \in \mathbb{R}$. Then $g h g^{-1} h^{-1}: z \mapsto z+(\lambda-1) b$. If $b \neq 0$, then this is a parabolic element that contradicts the assumption. If $b=0$, then $h: z \mapsto a z$. Thus, all elements of $G$ have the form $z \mapsto k z, k \in \mathbb{R}_{+}$. Then, up to an isomorphism, $G$ is a discrete subgroup of $\left(\mathbb{R}_{+}, \cdot\right)$. Then $G \cong \mathbb{Z}$.

Case 2. Let $|\mathcal{O}|=2$, say $\mathcal{O}=\left\{\alpha_{1}, \alpha_{2}\right\}$.
Fall 2.1. $\mathcal{O} \subseteq \mathbb{H}$.
Then all elements of $G \backslash\{e\}$ are elliptic. By Corollary 2.7.4, $G$ is a finite cyclic group.

Fall 2.2. $\mathcal{O} \subseteq \mathbb{R} \cup\{\infty\}$.
Then $G$ does not have parabolic elements. (Indeed, a parabolic element has only one finite orbit in $\mathbb{R} \cup\{\infty\}$, and this orbit contains only one point.) We consider three cases:
(a) $G \backslash\{e\}$ contains only hyperbolic elements. Then $\alpha_{1}$ and $\alpha_{2}$ are their common fixed points and both, $\left\{\alpha_{1}\right\}$ and $\left\{\alpha_{2}\right\}$ are orbits of $G$. Hence, $\mathcal{O}$ is not an orbit of $G$. A contradiction.
(b) $G \backslash\{e\}$ contains only elliptic elements. Then, by Corollary 2.7.4, $G$ is a finite cyclic group.
(c) $G \backslash\{e\}$ contains both, elliptic and hyperbolic elements. Then all elliptic elements have the order 2 and permute $\alpha_{1}, \alpha_{2}$. After an appropriate conjugation, we may assume that $\alpha_{1}=0$ and $\alpha_{2}=\infty$. Then all hyperbolic elements in $G$ have the form

$$
g_{k}: z \mapsto k z, \quad k>0, k \neq 1,
$$

and all elliptic elements in $G$ have the form

$$
e_{\lambda}: z \mapsto-\frac{\lambda}{z}, \quad \lambda>0 .
$$

Let $G_{0}$ be the subgroup of $G$ consisting of id and all hyperbolic elements of $G$. Then $G_{0}$ has index 2 in $G$. The second coset of $G_{0}$ in $G$ consists of elliptic elements. Let $e_{\lambda}$ be one of them. We conjugate $G$ by

$$
q: z \mapsto \sqrt{\lambda} z .
$$

Then $q G_{0} q^{-1}=G_{0}$ and $q e_{\lambda} q^{-1}=e_{1}=\psi$. Therefore

$$
q G q^{-1}=G_{0} \cup \psi G_{0} .
$$

Since $G_{0}$ is discrete, $G_{0}=\left\langle g_{k}\right\rangle$ for some $k>0, k \neq 1$. Hence $q G q^{-1}=H_{k}$.

Case 3. Let $|\mathcal{O}| \geqslant 3$.
Then $\mathcal{O} \subseteq \mathbb{R} \cup\{\infty\}$ and $G \backslash\{e\}$ contains only elliptic elements. By Corollary 2.7.4, $G$ is a finite cyclic group.

### 2.8 Jorgensen inequality

Lemma 2.8.1 Let $S, T \in \operatorname{PSL}_{2}(\mathbb{R})$. We set $S_{0}:=S$, and $S_{r+1}=S_{r} T S_{r}^{-1}$ for $r \geqslant 0$. If there is $n \geqslant 1$ such that $S_{n}=T$, then $\langle S, T\rangle$ is an elementary group and $S_{2}=T$.

Beweis. For $T=\mathrm{Id}$ the statement is evident. Let $T \neq \mathrm{Id}$.
Case 1. Let $|\widehat{\operatorname{Fix}}(T)|=1$.
Then $\widehat{\operatorname{Fix}}(T)=\{\alpha\}$ for some $\alpha \in \widehat{\mathbb{H}}$.
Since $S_{r+1}$ is conjugate to $T$, we have

$$
\begin{equation*}
\left|\widehat{\operatorname{Fix}}\left(S_{r+1}\right)\right|=1 \text { for all } r \geqslant 0 \tag{2.8.1}
\end{equation*}
$$

Moreover,

$$
S_{r+1} \circ S_{r}(\alpha)=S_{r} \circ T \circ S_{r}^{-1} \circ S_{r}(\alpha)=S_{r}(\alpha) .
$$

This implies

$$
\begin{equation*}
\widehat{\operatorname{Fix}}\left(S_{r+1}\right)=\left\{S_{r}(\alpha)\right\} . \tag{2.8.2}
\end{equation*}
$$

By assumption $S_{n}=T$, therefore $\widehat{\operatorname{Fix}}\left(S_{n}\right)=\widehat{\operatorname{Fix}}(T)=\{\alpha\}$. With the help of (2.8.1) and (2.8.2), we consequently obtain

$$
\begin{array}{ll}
\widehat{\operatorname{Fix}}\left(S_{n-1}\right) & =\{\alpha\}, \\
\ldots \\
\widehat{\operatorname{Fix}}\left(S_{1}\right) & =\{\alpha\}, \\
\widehat{\operatorname{Fix}}(S) & \supseteq\{\alpha\} .
\end{array}
$$

Therefore we have

1) $\alpha \in \widehat{\operatorname{Fix}}\langle S, T\rangle$. Hence, $\langle S, T\rangle$ is elementary.
2) $\widehat{\operatorname{Fix}}\left(S_{1}\right)=\{\alpha\}=\widehat{\operatorname{Fix}}(T)$. Hence, $S_{1} T=T S_{1}$ (see Theorem 2.6.2), i.e. $S_{2}=T$.

Case 2. Let $|\widehat{\operatorname{Fix}}(T)|=2$.
Then $\widehat{\operatorname{Fix}}(T)=\{\alpha, \beta\}$ for some $\alpha, \beta \in \widehat{\mathbb{H}}$. Since $S_{r+1}$ is conjugate to $T$, we have

$$
\begin{equation*}
\left|\widehat{\operatorname{Fix}}\left(S_{r+1}\right)\right|=2 \text {, and } S_{r+1} \text { is hyperbolic for all } r \geqslant 0 \tag{2.8.3}
\end{equation*}
$$

Moreover,

$$
S_{r+1} \circ S_{r}(\{\alpha, \beta\})=S_{r} \circ T \circ S_{r}^{-1} \circ S_{r}(\{\alpha, \beta\})=S_{r}(\{\alpha, \beta\}) .
$$

Each hyperbolic element has a unique invariant 2-elements set and this set coincides with the foxed point set. Therefore

$$
\begin{equation*}
\widehat{\operatorname{Fix}}\left(S_{r+1}\right)=\left\{S_{r}(\alpha, \beta)\right\} \tag{2.8.4}
\end{equation*}
$$

By assumption $S_{n}=T$, therefore $\widehat{\operatorname{Fix}}\left(S_{n}\right)=\widehat{\operatorname{Fix}}(T)=\{\alpha\}$. With the help of (2.8.1) and (2.8.2), we consequently obtain

$$
\begin{array}{ll}
\widehat{\operatorname{Fix}}\left(S_{n-1}\right) & =\{\alpha, \beta\}, \\
\ldots \\
\widehat{\operatorname{Fix}}\left(S_{1}\right) & =\{\alpha, \beta\}, \\
\widehat{\operatorname{Fix}}(S) & =\{\alpha, \beta\} .
\end{array}
$$

Therefore we have

1) $\{\alpha, \beta\} \supseteq \widehat{\operatorname{Fix}}\langle S, T\rangle$. Hence, $\langle S, T\rangle$ is elementary.
2) $\widehat{\operatorname{Fix}}\left(S_{1}\right)=\{\alpha, \beta\}=\widehat{\operatorname{Fix}}\left(S_{2}\right)$. Hence, $S_{1} S_{2}=S_{2} S_{1}$ (see Theorem 2.6.2) that implies $S_{1} T=T S_{1}$, i.e. $S_{2}=T$.

Theorem 2.8.2 (Jorgensen inequality) Let $T, S \in \mathrm{PSL}_{2}(\mathbb{R})$ such that $\langle T, S\rangle$ is a nonelementary Fuchsian group. Then

$$
\begin{equation*}
\left|\operatorname{Tr}^{2}(T)-4\right|+\left|\operatorname{Tr}\left(T S T^{-1} S^{-1}\right)-2\right| \geqslant 1 \tag{2.8.5}
\end{equation*}
$$

Proof. We set $S_{0}:=S$, and $S_{r+1}=S_{r} T S_{r}^{-1}$ for $r \geqslant 0$ as in Lemma 2.8.1. We will show that if the inequality (2.8.5) is not valid, then $S_{n}=T$ for some $n \in \mathbb{N}$. By Lemma 2.8.1 this would imply that the group $\langle S, T\rangle$ is elementary, a contradiction.

Case 1. Suppose that $T$ is parabolic.
Since the trace is invariant with respect to conjugation of matrices, we may assume that

$$
T=\left[\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right], \quad S=\left[\begin{array}{ll}
a & b \\
c & d
\end{array}\right]
$$

Suppose that the inequality (2.8.5) is not valid. A straightforward calculation shows that $|c|<1$. Let

$$
S_{r}=\left[\begin{array}{cc}
a_{r} & b_{r} \\
c_{r} & d_{r}
\end{array}\right]
$$

From $S_{r+1}=S_{r} \circ T \circ S_{r}^{-1}$ we obtain

$$
\left[\begin{array}{cc}
a_{r+1} & b_{r+1} \\
c_{r+1} & d_{r+1}
\end{array}\right]=\left[\begin{array}{cc}
a_{r} & b_{r} \\
c_{r} & d_{r}
\end{array}\right] \cdot\left[\begin{array}{cc}
1 & 1 \\
0 & 1
\end{array}\right] \cdot\left[\begin{array}{cc}
d_{r} & -b_{r} \\
-c_{r} & a_{r}
\end{array}\right]=\left[\begin{array}{cc}
1-a_{r} c_{r} & a_{r}^{2} \\
-c_{r}^{2} & 1+a_{r} c_{r}
\end{array}\right]
$$

By induktion we obtain $c_{r}=-c^{2^{r}}$. Since $|c|<1$, we have

$$
c_{r} \rightarrow 0
$$

The equality $a_{r+1}=1-a_{r} c_{r}$ and $|c|<1$ imply that

$$
\left|a_{r+1}\right| \leqslant 1+\left|a_{r} c_{r}\right| \leqslant 1+\left|a_{r}\right| \leqslant \cdots \leqslant(r+1)+|a| .
$$

Therefore $\left|a_{r} c_{r}\right| \leqslant(r+|a|)\left|c_{r}\right| \leqslant(r+|a|)|c|^{2^{r}} \rightarrow 0$. It follows

$$
a_{r+1}=1-a_{r} c_{r} \rightarrow 1
$$

Thus, $S_{r+1} \rightarrow T$. Since $\langle S, T\rangle$ is discrete, there exists $n$ such that $S_{n}=T$. Then by Lemma 2.8.1, the group $\langle S, T\rangle$ is elementary. A contradiction.

Case 2. Suppose that $T$ is hyperbolic.
After an appropriate conjugation, we may assume that

$$
T=\left[\begin{array}{cc}
u & 0 \\
0 & 1 / u
\end{array}\right], \quad S=\left[\begin{array}{ll}
a & b \\
c & d
\end{array}\right]
$$

We define

$$
\mu:=\left|\operatorname{Tr}^{2}(T)-4\right|+\left|\operatorname{Tr}\left(T S T^{-1} S^{-1}\right)-2\right| .
$$

Suppose that the inequality (2.8.5) is not valid. Then

$$
\mu=(1+|b c|)\left|u-\frac{1}{u}\right|^{2}<1 .
$$

In particular,

$$
\begin{equation*}
\left|\frac{1}{u}-u\right| \leqslant \mu^{\frac{1}{2}} \tag{2.8.6}
\end{equation*}
$$

From $S_{r+1}=S_{r} \circ T \circ S_{r}^{-1}$ follows

$$
\left[\begin{array}{ll}
a_{r+1} & b_{r+1} \\
c_{r+1} & d_{r+1}
\end{array}\right]=\left[\begin{array}{ll}
a_{r} d_{r} u-\frac{b_{r} c_{r}}{u} & a_{r} b_{r}\left(\frac{1}{u}-u\right) \\
c_{r} d_{r}\left(u-\frac{1}{u}\right) & \frac{a_{r} d_{r}}{u}-b_{r} c_{r} u
\end{array}\right]
$$

Then

$$
b_{r+1} c_{r+1}=-b_{r} c_{r}\left(1+b_{r} c_{r}\right)\left(u-\frac{1}{u}\right)^{2}
$$

Claim. We have $\left|b_{r} c_{r}\right| \leqslant \mu^{r}|b c|$.
Proof. We proceed the inductive step (where we use $\mu<1$ ):

$$
\left|b_{r+1} c_{r+1}\right| \leqslant \mu^{r}|b c| \cdot\left(1+\mu^{r}|b c|\right) \cdot\left|u-\frac{1}{u}\right|^{2} \leqslant \mu^{r}|b c| \cdot(1+|b c|) \cdot\left|u-\frac{1}{u}\right|^{2}=\mu^{r+1}|b c| .
$$

This claim and $\mu<1$ imply the following statements.
(1) $b_{r+1} c_{r+1} \rightarrow 0$,
(2) $a_{r} d_{r}=1+b_{r} c_{r} \rightarrow 1$,
(3) $a_{r+1} \rightarrow u$ and $d_{r+1} \rightarrow \frac{1}{u}$.

Then for all sufficiently large $r$ we have $\frac{a_{r}}{|u|} \leqslant \mu^{-\frac{1}{3}}$ (recall that $\mu<1$ ). This and (2.8.6) imply

$$
\frac{\left|b_{r+1}\right|}{\left|u^{r+1}\right|}=\frac{\left|b_{r}\right|}{\left|u^{r}\right|} \cdot \frac{\left|a_{r}\right|}{|u|} \cdot \frac{\left|\frac{1}{u}-u\right|}{1} \leqslant \frac{\left|b_{r}\right|}{\left|u^{r}\right|} \mu^{-\frac{1}{3}} \mu^{\frac{1}{2}}=\frac{\left|b_{r}\right|}{\left|u^{r}\right|} \mu^{\frac{1}{6}} .
$$

Since $\mu<1$, we deduce
(4) $\frac{\left|b_{r}\right|}{\left|u^{r}\right|} \rightarrow 0$. Analogously, we deduce
(5) $\left|c_{r}\right| \cdot\left|u^{r}\right| \rightarrow 0$.

Then

$$
T^{-r} S_{2 r} T^{r}=\left[\begin{array}{cc}
a_{2 r} & \frac{b_{2 r}}{u^{2 r}} \\
c_{2 r} u^{2 r} & d_{2 r}
\end{array}\right] \rightarrow T
$$

Since $\langle S, T\rangle$ is discrete, $T^{-r} S_{2 r} T^{r}=T$ for all sufficiently large $r$. Then $S_{2 r}=T$ for all sufficiently large $r$. We get a contradiction as in Case 1.

Case 3. Suppose that $T$ is elliptic.
By classification Lemma 2.1.2, the eigenvalues $\lambda_{1}, \lambda_{2}$ of $T$ satisfy $\lambda_{1}, \lambda_{2} \in \mathbb{C} \backslash \mathbb{R}$, $\lambda_{1}=\bar{\lambda}_{2}, \lambda_{1} \lambda_{2}=1$. Hence, $\lambda_{1}=e^{i \varphi}$ for some $\varphi \in(0, \pi) \cup(\pi, 2 \pi)$. Thus, after an appropriate conjugation in $\mathrm{PSL}_{2}(\mathbb{C})$, we may assume that

$$
T=\left[\begin{array}{cc}
u & 0 \\
0 & 1 / u
\end{array}\right] \quad\left(u=e^{i \varphi}, 0<\varphi<\pi\right) .
$$

The remaining proof follows repeats that in Case 2.


[^0]:    ${ }^{1}$ This will be clear only after Theorem 1.4.4.

[^1]:    ${ }^{2}$ In this case, we have $f:(x, y) \mapsto \frac{1}{y^{2}}$ and $\varphi:(x, y) \mapsto(u, v)$.

