# Scuola Normale Superiore di Pisa 

 Classe di Scienze
## Wilhelm Klingenberg <br> Closed geodesics on surfaces of genus 0

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# Closed Geodesics on Surfaces of Genus $0{ }^{(*)}$. 

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1.     - One of the outstanding results in Global Differential Geometry is the theorem of Lusternik and Schnirelmann that there exist on any closed orientable surface of genus 0 three non-selfintersecting closed geodesics. See [LS 1, 2] as well as our «Lectures on Closed Geodesics» [Kl]-briefly: LCG-for a simplified proof.

In the present paper we want to show how an extension of the methods which we employed in LCG for the proof of the Lusternik-Schnirelmann theorem yields an elementary proof for the existence of infinitely many prime (i.e., not multiply covered) closed geodesics on a surface of genus 0 .

In LCG we have already presented a proof for the existence of infinitely many prime closed geodesics on an arbitrary riemannian manifold with finite fundamental group. This proof, however, is quite involved and uses rather delicate arguments. So it seems of interest to give a much simpler proof for the case of a surface of genus 0 .

A fundamental rôle in our proof play the subspaces of the space of all closed curves formed by those curves having selfintersection number $\leqslant \mathbf{a}$ fixed integer $\nu$. We will introduce energy-decreasing deformations of those spaces into itself and show that for every $\nu=0,1, \ldots$, there exist two not necessarily prime closed geodesics with selfintersection number $\nu$. Moreover, as $\nu$ increases, the energy of these geodesics does. Cf. (5.5) for details.
2. - We begin by fixing some notations. By $M$ we denote the 2 -sphe re $S^{2}$, endowed with an arbitrary riemannian metric. $P M$ shall be the space of all
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piecewise differentiable maps (with an infinite number of geodesic pieces being allowed)

$$
c: S=[0,1] /\{0,1\} \rightarrow M
$$

of the parameterized circle $S$ into $M$, endowed with the metric

$$
d_{\infty}\left(c, c^{\prime}\right)=\sup _{t \in S} d\left(c(t), c^{\prime}(t)\right)
$$

Here $d($,$) denotes the distance on M$, derived from the riemannian metric.
On $P M$ we have defined the energy

$$
E: P M \rightarrow \boldsymbol{R} ; \quad c \mapsto \frac{1}{2} \int_{S}|\dot{c}(t)|^{2} d t
$$

For any real $x$ we put

$$
P^{\kappa} M=\{c \in P M ; E(c) \leqslant \varkappa\} .
$$

In particular, $P^{0} M$ consists of the constant maps $c: S \rightarrow M$ and thus is canonically isomorphic to $M$.

On $P M$ we have the canonical $S$-action coming from the $S=S O(2)$ action on $S$ :

$$
(z, c)=\left(e^{2 \pi i r}, c(t)\right) \in S \times P M_{\mapsto} \cdot z \cdot c=(c(t+r)) \in P M .
$$

That is to say, $z \cdot e$ is obtained from $c$ by changing the initial point from $0 \in S$ to $r \in S$.

We also have the $\boldsymbol{Z}_{2}$-action generated by

$$
\theta: P M \rightarrow P M ; \quad(c(t)) \mapsto(c(1-t))
$$

i.e., by the reversal of the orientation.

The $S$-action and the $Z_{2}$-action together give a $O(2)$-action on $P M$.
The isotropy group $I^{\sim}(c)$ of $c \in P M$ under the $S$-action either is $S$ (and this occurs if and only if $c \in P^{0} M$ ), or else, is a finite cyclic group. The order of this group is called the multiplicity of $c$. If $c$ has multiplicity 1 it also is called prime. If $c$ has multiplicity $m$ then $c$ can be written as $c(t)=c_{0}(m t) . \quad c_{0}$ is uniquely determined and is called the underlying prime closed geodesic of $c$.

We now introduce the fundamental concept of the selfintersection number $\nu(c)$ of a closed curve $c \in P M$. Assume first that $c: S \rightarrow M$ is an immersion and that the multiple points-if they occur at all-all are double points
representing a transversal intersection. Then we define $\nu(c)$, to be the number of these double points.

For every integer $\nu \geqslant 0$ we now define $P_{\nu} M$ to be the closure in $P M$ of the immersions with $\leqslant \nu$ transversal double points. Put $P_{-1} M=\emptyset$. We thus have the strictly increasing sequence

$$
P_{0} M \subset P_{1} M \subset P_{2} M \subset \ldots
$$

of subspaces of $P M$. If $c \in P_{\nu} M$ but not in $P_{\nu-1} M$ then we define $\nu(c)$ to be $\nu$.

Not every $c \in P M$ will belong to one of the $P_{\nu} M$. For the remaining $c$ we put $\nu(c)=\infty$. But certainly for closed geodesics $c, v(c)$ is finite.

Actually, there exists the following relation between the selfintersection number of a multiply covered closed geodesic and the selfintersection number of the underlying prime closed geodesic:
2.1 Proposition. Let $c$ be a closed geodesic of multiplicity $m(c)$. Denote by $c_{0}$ the underlying prime closed geodesic. Then

$$
\nu(c)=m(c)^{2} v\left(c_{0}\right)+m(c)-1 .
$$

Proof. We can asume: $m(c)>1$. Observe now that the $m(c)$-fold covered closed geodesic $c_{0}$ can be approximated by curves which have $m(c)$ arcs near $c_{0}$ and $m(c)-1$ transversal double points away from the selfintersection points of $c_{0}$. Near each of the selfintersection points of $c_{0}$ (which are transversal but may not be only double points but also points with $k>2$ arcs passing through it) the approximating curves can be assumed to have double points only in number $m(c)^{2} k$ if $k+1$ is the number of arcs passing through that selfintersection point.

Note. It is the concept of the selfintersection number where it is important that $M$ is 2 -dimensional. Indeed, if $\operatorname{dim} M>2, P_{0} M=P M$ and thus the selfintersection number carries no information in this case.
3. - Our next goal is to define an $E$-decreasing deformation of the subspace $P_{\nu} M$ into itself. For $\nu=0$, Lusternik and Schnirelmann, cf. [Ly], had defined already such a deformation. This deformation seems very complicated and we had difficulties in establishing the necessary properties for it. Therefore we proposed in LCG a much simpler and actually more efficient $E$-decreasing deformation of $P_{0} M$ into itself. Here we will show that the same deformation also can be used to transform $P_{\nu} M$ into itself.

As it is customary in this theory, the deformations are defined only on a subspace $P_{\nu}^{\chi} M=P_{v} M \cap P^{\chi} M$ of $P_{\nu} M, x$ some arbitrarily fixed positive real number. For such a $x$ we choose an even integer $k>0$ such that $4 \varkappa / k \leqslant \eta^{2}$. Here, $\eta$ is a positive number such that, for every $p \in M$, the disc of radius $2 \eta$ around $p$ is a convex neighborhood of $p$.

Let $c \in P M$. The relation (with $L$ denoting the length)

$$
L\left(c \mid\left[t_{0}, t_{1}\right]\right)^{2} \leqslant 2 E\left(c \mid\left[t_{0}, t_{1}\right]\right)\left|t_{1}-t_{0}\right|
$$

implies: If $c \in P^{\kappa} M$ and $\left|t_{1}-t_{0}\right| \leqslant 2 / k$ then

$$
L\left(c \mid\left[t_{0}, t_{1}\right]\right) \leqslant \eta .
$$

Thus in particular, there exist a unique minimizing geodesic segment $c_{c\left(t_{0}\right) c\left(t_{1}\right)}$ from $c\left(t_{0}\right)$ to $c\left(t_{1}\right)$.

We begin by defining the deformation

$$
\tilde{\mathfrak{D}}_{\sigma}: P_{\nu}^{\kappa} M \rightarrow P_{\nu}^{\kappa} M ; \quad \sigma \in[0,1 / k] .
$$

Let $c \in P_{\nu}^{*} M$. Put $c(0)=p$. For every $\sigma \in[0,2 / k]$ the geodesic segment $c_{p c(\sigma)}$ is well defined. The closed curve $\mathscr{D}_{\sigma} c=c_{p c(\sigma)} \cup c \mid[\sigma, 1]$ again belongs to $P^{\kappa} M$. However, it may happen that, as $\sigma$ increases, $c_{p c(\sigma)}$ will have new proper (i.e., transversal) intersections with $c \mid[\sigma, 1]$ as compared to the original arc $c \mid[0, \sigma]$. Whenever this begins to occur we start modifying $\sigma \mid[\sigma, 1]$ by substituting for each small are of $c \mid[\sigma, 1]$, which comes to lie on the «wrong» side of $c_{p c(\sigma)}^{\prime}$ so as to cause additional transversal selfinteractions, a geodesic segment on $c_{p c(\sigma)}$ which goes from the initial point to the end point of such an arc.

It may be necessary to make such a substitution simultaneously for several small arcs. But this is no problem, since the whole procedure takes place inside the convex neighborhood $B_{2 \eta}(p)$ of $p$ where the geodesic segments $c_{p c(\sigma)}$ look like straight segments starting from $p$. The substituting process may be thought of as a pushing aside or sweeping aside of parts of $c \mid[\sigma, 1]$ which come to lie in the way of the segment $c_{p c(\sigma)}$ as it moves around.

We denote by $\widetilde{\mathscr{D}}_{\sigma} c \mid[\sigma, 1]$ the modification of $c \mid[\sigma, 1]$. Since arcs of $c \mid[\sigma, 1]$ may be replaced by geodesic segments, the $E$-value of $\tilde{\mathscr{D}}_{\sigma} c \mid[\sigma, 1]$ will be $\leqslant E(c \mid[\sigma, 1])$. As we proceed with $\sigma$ from 0 to $2 / k$ it also may happen that $c(\sigma)$ has been replaced by $\widetilde{\mathfrak{D}}_{\sigma} c(\sigma)$. We now define

$$
\tilde{\mathscr{D}}_{\sigma} c=c_{p} \tilde{\mathscr{D}}_{o} c(\sigma) \cup \widetilde{\mathscr{D}}_{\sigma} c \mid[\sigma, 1]
$$

We proceed to define in the same way $\widetilde{\mathscr{D}}_{\sigma} c$ for $\sigma \in[2 / k, 4 / k], c \in P_{\nu}^{\alpha} M$ : The only difference is that $c(0)$ is being replaced by $c(2 / k)$ as initial point and we begin by substituting the geodesic segment $c_{c(2 / k) c(\sigma)}$ for the arc $c\left|[2 / k, \sigma] . \widetilde{D}_{4 / k} c\right|[2 / k, 4 / k]$ will be a geodesic segment of length $\leqslant \eta$.

We continue in this manner until we reach the interval $[(k-2) / k, 1]$. We then proceed to define $\widetilde{D}_{\sigma} c$ for $\sigma \in[1,1+2 / k]$ as being a similar deformation with $c(1 / k)$ as initial point and having the effect that $c \mid[1 / k, \sigma-1+1 / k]$, $\sigma \in[1,1+2 / k]$ is being replaced by a geodesic segment. We go on until we reach the interval $\sigma \in[2-2 / k, 2]$ where $\widetilde{\mathscr{D}}_{\sigma} c$ represents a curve such that $\widetilde{D}_{\sigma} c \mid[1-1 / k, \sigma-1+1 / k]$ is a geodesic segment.

We now define $\widetilde{\mathscr{D}}(\sigma, c), \sigma \in[0,2]$, to be the subsequent application of the mappings $\widetilde{\mathscr{D}}_{2 / k}, \ldots, \widetilde{\mathfrak{D}}_{2 l / k}, \widetilde{\mathfrak{D}}_{\sigma}$, where $2 l$ is the even integer determined by $2 l \leqslant k \sigma<2 l+2$.
3.1 Proposition. Choose $a x>0$ and an integer $\nu \geqslant 0$. Then the mapping

$$
\begin{gathered}
\widetilde{D}:[0,2] \times P_{\nu}^{\alpha} M \rightarrow P_{\nu}^{\alpha} M \\
(\sigma, c) \mapsto \mathscr{D}(\sigma, c)
\end{gathered}
$$

is continuous. Moreover, $E(\widetilde{D}(2, c)) \leqslant E(c)$ with equality if and only if either $c$ is a constant map or else a closed geodesic.

Proof. We consider $P_{\nu}^{\alpha} M$ to be a subspace of $P M$. Then the continuity is obvious. The last statement follows from the following standard facts of riemannian geometry: Let $c:\left[t_{0}, t_{1}\right] \rightarrow M$ be a piecewise differentiable $\operatorname{map}$ from $p=c\left(t_{0}\right)$ to $q=c\left(t_{1}\right), d(p, q) \leqslant \eta$. Let $c_{p q}$ be the unique minimizing geodesic segment from $p$ to $q$, parameterized by $\left[t_{0}, t_{1}\right]$. Then $E\left(c_{p q}\right) \leqslant E(c)$ with equality if and only if $c_{p q}(t)=c(t)$.

For reference we formulate another standard result of local riemannian geometry:
3.2 Proposition. Let $\left\{c_{n}\right\}$ be a sequence of piecewise differentiable paths $c_{n}:[0,1] \rightarrow$ M. Put $c_{n}(0)=p_{n}, c_{n}(1)=q_{n}$. Assume: $d\left(p_{n}, q_{n}\right) \leqslant \eta$ with $\eta>0$ as above so that there exists the uniquely determined minimizing geodesic segment $c_{p_{n} q_{n}}:[0,1] \rightarrow M$ from $p_{n}$ to $q_{n}$.

Let now $\left\{E\left(c_{n}\right)\right\}$ and $\left\{E\left({c_{p_{n} \sigma_{n}}}\right)\right\}$ both be convergent with the same limit. Then $\left\{c_{n}\right\}$ possesses a cunvergent subsequence with limit a minimizing geodesic segment $c:[0,1] \rightarrow M$.

Proof. Since $M$ is compact, there exists a subsequence of $\left\{c_{n}\right\}$ which we denote again by $\left\{c_{n}\right\}$ such that $\lim p_{n}=p$ and $\lim q_{n}=q$ exist. It fol-
lows that $\left\{c_{v_{n} q_{n}}\right\}$ converges to the minimizing geodesic segment $c=c_{p q}$ from $p$ to $q$.

Let $\left.t_{0} \in\right] 0,1\left[\right.$. Put $c_{n}\left(t_{0}\right)=r_{n}$. We will show that $\lim r_{n}$ exists and is equal to $c\left(t_{0}\right)$. Indeed, $\left\{r_{n}\right\}$ possesses a convergent subsequence $\left\{r_{n(k)}\right\}$. Let $r$ be its limit. The sequences

$$
\left\{c_{p_{n(k)} r_{n(k)}}\right\}, \quad\left\{c_{\left.r_{n(k)}\right)_{n(k)}}\right\}
$$

of unique minimizing geodesic segments converge to $c_{p r}$ and $c_{r q}$, respectively, with parameter domain $\left[0, t_{0}\right]$ and $\left[t_{0}, 1\right]$. From our assumptions follows:

$$
E\left(c_{p r}\right)+E\left(c_{r q}\right) \leqslant E\left(c_{p q}\right) .
$$

Since $c_{p q}:[0,1] \rightarrow M$ is the unique segment of minimal $E$-value from $p$ to $q$, the curve $c_{p r} \cup c_{r q}$ from $p$ to $q$ must coincide with $c_{p q}$. In particular, its value $r$ for $t=t_{0}$ must be equal to $c_{p q}\left(t_{0}\right)$. This completes the proof.
3.3 Lemma. Choose $x \geqslant 0$ and an integer $v \geqslant 0$. Let $\widetilde{D}(\sigma$,$) be the deforma-$ tion of $P_{\nu}^{*} M$ into itself, as defined above.

Assume that $\left\{c_{n} \nu\right\}$ is a sequence in $P_{\nu}^{\alpha} M$ such that both, $\left\{E\left(c_{n}\right)\right\}$ and $\left\{E\left(\widetilde{D}\left(2, c_{n}\right)\right)\right\}$, are convergent with the same limit $\varkappa_{0}>0$. Then $\left\{c_{n}\right\}$ possesses a convergent subsequence with limit a closed geodesic $c_{0}$ with $E\left(c_{0}\right)=\varkappa_{0}, \nu\left(c_{0}\right) \leqslant \nu$. If, for all sufficiently large, $\nu\left(\widetilde{D}\left(2, c_{n}\right)\right)=\nu$ then also $\nu\left(c_{0}\right)=\nu$.

Proof. Besides the deformation $\widetilde{D}(\sigma$,$) we will also consider the defor-$ mation $\mathfrak{D}(\sigma$, ) of Lusternik and Schnirelmann, (cf. [Ly] and LCG) which is defined as follows: Choose $k$ as for the definition of $\widetilde{D}(\sigma$,$) .$

For $\sigma \in[j / k,(j+2) / k], j=0,2, \ldots, k-2$, put

$$
\begin{aligned}
& D_{\sigma} c \mid[j / k, \sigma]=c_{c(j / k) c(\sigma)} \\
& \mathscr{D}_{\sigma} c|[\sigma, j / k+1]=c|[\sigma, j / k+1]
\end{aligned}
$$

Similarly, for $\sigma \in[1+j / k, 1+(j+2) / k], j=0,2, \ldots, k-2$, define

$$
\begin{aligned}
& \mathfrak{D}_{\sigma} c \mid[(j+1) / k, \sigma-1+1 / k]=c_{c((j+1) / k) c(\sigma-1+1 / k)} \\
& \mathscr{D}_{\sigma} c|[\sigma-1+1 / k, 1+(j+1) / k]=c|[\sigma-1+1 / k, 1+(j+1) / k]
\end{aligned}
$$

Here, as always, the $t$-parameter has to be taken modulo 1.
Let now $\mathfrak{D}(\sigma, c), \sigma \in[0,2]$, be the subsequent application of the deformations $\mathfrak{D}_{2 / k}, \ldots, \mathfrak{D}_{2 l / k}, \mathscr{D}_{\sigma}$, with $2 l$ being the even integer determined by $2 l \leqslant k \sigma<2 l+2$.

We then have the relation

$$
E\left(\widetilde{D}\left(\sigma, c_{n}\right)\right) \leqslant E\left(\mathscr{D}\left(\sigma, c_{n}\right)\right) \leqslant E\left(c_{n}\right) .
$$

Thus we get from our hypothesis and from (3.2), since $\mathfrak{D}\left(1, c_{n}\right)$ consists of $k / 2$ geodesic segments: There exists a subsequence of $\left\{c_{n}\right\}$ which we denote again by $\left\{c_{n}\right\}$ such that both, $\left\{c_{n}\right\}$ and $\left\{D\left(1, c_{n}\right)\right\}$ converge to a $k / 2$ times broken geodesic $c_{0}$ with $E\left(c_{0}\right)=\varkappa_{0}$. Applying the same argument to the sequences

$$
\left\{\mathscr{D}\left(1, \mathfrak{D}\left(1, c_{n}\right)\right)\right\}=\left\{\mathscr{D}\left(2, c_{n}\right)\right\} \quad \text { and } \quad\left\{\mathscr{D}\left(1, c_{n}\right)\right\}
$$

we see: $\left\{\mathscr{D}\left(2, c_{n}\right)\right\}$ converges to $c_{0}$. Since $E\left(\widetilde{D}\left(2, c_{0}\right)\right)=E\left(\mathscr{D}\left(2, c_{0}\right)\right)=E\left(c_{0}\right)$, $c_{0}$ is a closed geodesic, cf. (3.1).

Clearly $\nu\left(c_{0}\right) \leqslant \nu$, since $P_{v}^{*}$ is closed. If $\nu\left(\mathscr{D}\left(2, c_{n}\right)\right)=v$, all large $n$, then also $\nu\left(c_{0}\right)=\nu$. To see this we observe that $\mathfrak{D}\left(2, c_{n}\right)$ consists of $k / 2$ geodesic segments: $\mathfrak{D}\left(2, c_{n}\right) \mid[(j+1) / k,(j+3) / k]$ is such a segment and it converges to $c_{0} \mid[(j+1) / k,(j+3) / k]$. Hence, $v\left(D\left(2, c_{n}\right)\right)=v\left(c_{0}\right)$. But

$$
\nu \geqslant v\left(c_{0}\right)=v\left(\mathfrak{D}\left(2, c_{n}\right)\right) \geqslant v\left(\widetilde{D}\left(2, c_{n}\right)\right)=v
$$

hence our claim. This completes the proof of (3.3).
From (3.3) we easily get
3.4 Lemma. Let $\chi_{0}>0$. Let $v$ be an integer $\geqslant 0$. Denote by $C$ the set of closed geodesics $e$ with $E(c)=\varkappa_{0}, \nu(c) \leqslant \nu$. Let $ひ$ be an open neighborhood of $C$. In the case $C=\emptyset$ one may choose $\mathcal{U}=\emptyset$. Let $\varkappa>\varkappa_{0}$ and consider on $P_{\nu}^{*} M$ the deformation $\widetilde{D}=\widetilde{D}(2$,$) . Then there exists a \varepsilon>0$ such that

$$
\widetilde{D} P_{\nu}^{x_{0}+\varepsilon} \subset \mathcal{U} \cup P_{\nu}^{x_{0}-\varepsilon} .
$$

Proof. Since $\widetilde{D}$ is continuous, there exists an open neighborhood $\mathcal{U}^{\prime}$ of $C$, U' $^{\prime} \mathcal{U}$, such that $\widetilde{D} \omega^{\prime} \subset \mathcal{U}$. If there were no $\varepsilon>0$ with the desired property this would imply the existence of a sequence $\left\{c_{n}\right\}$ in $P_{\nu}^{\alpha} M$ with $c_{n} \notin$ U' $^{\prime}$,

$$
\varkappa_{0}-1 / n \leqslant E\left(\widetilde{D} c_{n}\right) \leqslant E\left(c_{n}\right) \leqslant \varkappa_{0}+1 / n .
$$

From (3.3) we get a convergent subsequence of $\left\{c_{n}\right\}$ with limit a closed geodesic $c_{0}, E\left(c_{0}\right)=\varkappa_{0}, c_{0} \in P_{\nu}, c_{0} \notin U^{\prime}$. This clearly is impossible.
4. - In this paragraph we are going to construct for every $\nu=0,1, \ldots$, a rationally non-trivial cycle $w(v)$ of $P M$ which actually belongs to $P_{\boldsymbol{v}} M$.

In addition, we will define a certain $\boldsymbol{Z}_{2}$-cycle $w^{\prime}(\nu)$ of the quotient space $\boldsymbol{P}_{\boldsymbol{\nu}} M / \vartheta$. In (5) this will be used to prove the existence of two closed geodesics with selfintersection number $\nu$.

Topologically, i.e., if we disregard the distance but only consider the underlying topology, $P M$ can be identified with $P S^{2}$. Now, the homology of $P S^{2}$ is well known. An appropriate way to describe $P S^{2}$ is to consider it as subset of the Hilbert manifold $\Lambda S^{2}$ of all closed $H^{1}$-curves $c: S \rightarrow S^{2}$, cf. LCG for more details. $\mathbb{S}^{2}$ here is being endowed with the canonical metric of constant curvature equal to 1 .

The energy function $\widetilde{E}: \Lambda S^{2} \rightarrow \boldsymbol{R}$ is differentiable. It satisfies the condition (C) of Palais and Smale. The set of critical points of $\tilde{E}$ on $\Lambda S^{2}$ decomposes into nondegenerate critical submanifolds: For $\tilde{E}=0$, this is the manifold $\Lambda^{0} S^{2}$ of constant maps, isomorphic to $S^{2}$. For $\widetilde{E}=2(v+1)^{2} \pi^{2}$, $\nu=0,1,2, \ldots$, we have the manifold $B_{v+} S^{2}$ of the $(\nu+1)$-fold covered great circles. Each $B_{v+1} S^{2}$ is isometric to $B_{1} S^{2}$, i.e., to the unit tangent bundle of $S^{2}$, which coincides with the real projective space $P^{3}$. Index $B_{\nu+1} S^{2}=2 v+1$.

We therefore obtain $\Lambda^{2 \pi^{2}(v+1)^{2}} S^{2}$ by attaching to $\Lambda^{2 \pi^{2} r^{2}} S^{2}$ the negative bundle of $B_{v+1} S^{2}$. One can show that this attaching does not affect the homology of $\Lambda^{2 \pi^{2} \nu^{2}} S^{2}$. Since $H_{i}\left(B_{\nu+1} S^{2}\right) \neq 0$ only for $i=0$ and 3 , this attaching gives new rational cycles in dimensions $2 v+1$ and $2 v+4$.

For our purposes it seems preferable to describe the $(2 v+1)$-dimensional homology class determined by the negative fibre of an element $c_{v+1} \in B_{v+1} S^{2}$ in a different manner: We take the standard embedding of $S^{2}$ in $\boldsymbol{R}^{3}=\left\{\left(x_{0}, x_{1}, x_{2}\right)\right\}$. As base point $*$ of $S^{2}$ we choose $(1,0,0)$. Let

$$
a:\left(S^{2}, *\right) \rightarrow\left(S^{2}, *\right) ; \quad b:\left(S^{3}, *\right) \rightarrow\left(S^{2}, *\right)
$$

be the identity map and the Hopf map, respectively. To give a precise description of $b$ we represent $S^{3}$ by

$$
\left\{\left(w_{0}, w_{1}\right) \in \boldsymbol{C} \times \boldsymbol{C} ; w_{0} \bar{w}_{0}+w_{1} \bar{w}_{1}=1\right\}
$$

The base point $*$ of $S^{3}$ shall be $(1,0) \in C \times C$. Then

$$
b:\left(w_{0}, w_{1}\right) \in S^{3} \mapsto z=w_{0} / w_{1} \in \bar{C}=C \cup \infty \cong S^{2}
$$

Here, the isomorphism between $\overline{\boldsymbol{C}}$ and $S^{2}$ shall be the stereographic projection of $\mathbb{S}^{2}-\{*\}$ onto $\boldsymbol{C} \cong \boldsymbol{R}^{2}=\{(x, y)\}$ :
$x_{0}=\left(x^{2}+y^{2}-1\right) /\left(x^{2}+y^{2}+1\right) ; x_{1}=2 x /\left(x^{2}+y^{2}+1\right) ; x_{2}=2 y /\left(x^{2}+y^{2}+1\right)$.
$\infty \in \overline{\boldsymbol{C}}$ then corresponds to $* \in \boldsymbol{S}^{2}$.

We also describe $S^{3}$ by

$$
\left\{\left(e^{i \varphi} \cos \alpha, e^{-i \psi} \sin \alpha\right) ; 0 \leqslant \alpha \leqslant \pi ; 0 \leqslant \varphi, \psi \leqslant 2 \pi\right\}
$$

Then

$$
b\left(e^{i \varphi} \cos \alpha, e^{-i \varphi} \sin \alpha\right)=(\cos 2 \alpha, \sin 2 \alpha \cos (\varphi+\psi), \sin 2 \alpha \sin (\varphi+\psi))
$$

On $S^{3}$ we consider the 2 -sphere $S_{0}^{2}$ given by $\left\{w_{0}-\bar{w}_{0}=0\right\}$, i.e.

$$
S_{0}^{2}=\left\{\left(\cos \alpha, e^{-i \varphi} \sin \alpha\right) ; 0 \leqslant \alpha \leqslant \pi ; 0 \leqslant \psi \leqslant 2 \pi\right\}
$$

Every $q \in S_{0}^{2}$ determines a parameterized circle $c_{q}$ on $S^{3}$ as follows: $c_{q}$ starts from $* \in S_{0}^{2} \subset S^{3}$ tangentially to the great circle $\left\{\left|w_{0}\right|=1\right\}$ into the halfsphere $\left\{-i\left(w_{0}-\bar{w}_{0}\right) \geqslant 0\right\}$ and passes for $t=\frac{1}{2}$ through the point $q \in S_{0}^{2}$.

If $q$ is being given by $\left(\cos \alpha_{0}, e^{-i \varphi_{0}} \sin \alpha_{0}\right)$ then the circle $c_{q}$ is determined by the following two linear equations for $u_{0}, v_{0}, u_{1}, v_{1}$, with $w_{0}=u_{0}+i v_{0}$, $w_{1}=u_{1}+i v_{1}$ :
$u_{1} \sin \psi_{0}+v_{1} \cos \psi_{0}=0 ; \quad u_{0} \sin \alpha_{0}-u_{1} \cos \alpha_{0} \cos \psi_{0}+v_{1} \cos \alpha_{0} \sin \psi_{0}=0$.
The image under $b: S^{3} \rightarrow S^{2}$ of the circle $c_{q}$ therefore is the circle $\bar{b}(q)$, passing through $* \in \mathbb{S}^{2}$, which in $C$ is being represented by the straight line

$$
\cos \psi_{0} x+\sin \psi_{0} y=\operatorname{ctg} \alpha_{0}
$$

In particular, to $\alpha_{0}=\pi / 2$, i.e., the points $q$ on the great circle $\left\{\left(0, e^{-i \varphi}\right)\right\} \subset$ $\subset S_{0}^{2} \subset S^{3}$ there correspond the great circles on $S^{2}$ which pass through $* \in S^{2}$.

We have thus defined a mapping

$$
\bar{b}:\left(S_{0}^{2}, *\right) \rightarrow\left(\Omega S^{2}, *\right)
$$

where $\Omega S^{2}$ is the loop space of $S^{2}$. We consider only the piecewise differentiable loops starting and ending at $* \in S^{2}$. Thus, $\Omega S^{2}$ becomes a subspace of $P S^{2}$.

We also define a map

$$
\bar{a}:\left(S_{0}^{1}, *\right) \rightarrow\left(\Omega S^{2}, *\right)
$$

as follows: Let $S_{0}^{1}$ be the great circle on $S^{2}$ given by $\left\{x_{1}=0\right\}$. Then $\bar{a}(p)$ shall be the parameterized circle starting from $* \in S_{0}^{1} \subset S^{2}$ into the half sphere $\left\{x_{1} \geqslant 0\right\}$ tangentially to the great circle $\left\{x_{2}=0\right\}$ and passing for $t=\frac{1}{2}$ through the point $p \in S_{0}^{1}$.

The mappings $\bar{a}$ and $\bar{b}$ are but particular explicit descriptions of the mappings $\Omega a \in \pi_{1} \Omega S^{2}, \Omega b \in \pi_{2} \Omega S^{2}$, associated to $a \in \pi_{2} S^{2}, b \in \pi_{3} S^{2}$, cf. Spanier [Sp]. Thus, $\bar{a}, \bar{b}$ are cycles which represent the generators of the rational homology of $\Omega S^{2}$ in dimensions 1 and 2. A classical result of H. Hopf states that the rational cohomology ring of $\Omega S^{2}$ is being generated by the duals $[\bar{a}]^{*},[\bar{b}]^{*}$ of the homology classes $[\bar{a}]$ and $[\bar{b}]$. Moreover-and this is fundamental for our further argumentation-the products $[\bar{a}]^{*} \cup[\bar{b}]^{* v}, v=0,1, \ldots$, are restrictions of cohomology classes of $P S^{2} \cong \Lambda S^{2}$ (homotopy equivalence), cf. Švarc [Šv].

This shows that we can define cycles $w(\nu)$ of $P M$, dual to the $\left[\bar{a}^{*}\right] \cup\left[\bar{b}^{*}\right]^{\nu}$, by taking the loop product (also called: Pontrjagin product, cf. BottSamelson [BS]) of $\bar{a}$ with $v$ copies of $\bar{b}$ :

$$
w(v):(p, q) \in S_{0}^{1} \times\left(S_{0}^{2}\right)^{\nu} \mapsto \bar{a}(b) \cdot \bar{b}(q) \in \Omega S^{2}
$$

Here, $q$ stands for the $\nu$-tupel $\left(q_{1}, \ldots, q_{v}\right)$ and $\bar{b}(q)$ denotes the loop product $\bar{b}\left(q_{1}\right) \cdot \ldots \cdot \bar{b}\left(q_{v}\right)$. In forming the loop product of circles we take the parameterization of the factors so as to be proportionally to arc length. I.e.:

$$
L\left(w(v)(p, q) \mid\left[0, t_{0}\right]\right)=t_{0} L(w(v)(p, q))
$$

We also use $D(v)$ to denote the domain of $w(v)$. As we stated before, $w(v)$ is homologous to the negative fibre of the $(\nu+1)$-fold covering of a great circle on $S^{2}$.

By $\gamma_{\tau}$ we denote the positive rotation of $S_{0}^{2}$ around $* \in S_{0}^{2}$ and its antipodal point $*^{\prime} \in S_{0}^{2}$ by the angle $\tau$. I.e. if $q$ is being represented by $\left(\cos \alpha, e^{-i \varphi} \sin \alpha\right)$ then $\gamma_{\tau} q$ shall be given by $\left(\cos \alpha, e^{-i(\varphi+\tau)} \sin \alpha\right)$.

We denote by

$$
\iota: S_{0}^{2} \rightarrow S_{0}^{2} ; \quad\left(\cos \alpha, e^{-i \varphi} \sin \alpha\right) \mapsto\left(-\cos \alpha, e^{-i \varphi} \sin \alpha\right)
$$

the reflection on the great circle $\left\{\left|w_{1}\right|=1\right\} \subset S_{0}^{2}$. Then

$$
\lambda=\gamma_{\pi} \circ \iota=\iota \circ \gamma_{\pi}: S_{0}^{2} \rightarrow S_{0}^{2}
$$

is the antipodal map.
Finally, we use $\gamma_{\tau}$ also to denote the positive rotation of $\boldsymbol{S}^{2} \subset \boldsymbol{R}^{3}$ around the points $*=(1,0,0), *^{\prime}=(-1,0,0)$ by the angle $\tau$.

### 4.1 Proposition.

(i) $\gamma_{\tau} \bar{b}(q)=\bar{b}\left(\gamma_{\tau} q\right) ;$
(ii) $\bar{b}(\lambda q)=\theta \bar{b}(q)$;
(iii) $w(\nu+1)(*, \lambda q)=\theta w(v+1)(*, q)$,
with $\lambda q=\lambda\left(q_{1}, \ldots, q_{v}\right)=\left(\lambda q_{v}, \ldots, \lambda q_{1}\right)$.
Proof. (i) follows from $z\left(\gamma_{\tau} q\right)=e^{i \tau} z(q)$. To see (ii) we write $\bar{b}(\lambda q)=$ $=\gamma_{\pi} \bar{b}(\iota q)$. Let $q$ have $(\alpha, \psi)$ as coordinates. The linear equation of the straight line representing the circle $\bar{b}(q)$ in the plane $\boldsymbol{C}$ is then given by

$$
x \cos \psi+y \sin \psi=\operatorname{ctg} \alpha
$$

Therefore, $\bar{b}(\iota q)$ possesses the representation

$$
x \cos \psi+y \sin \psi=-\operatorname{ctg} \alpha
$$

which means that $\gamma_{\pi} \bar{b}(\iota q)$ coincides with $\bar{b}(q)$, up to the orientation.
To prove (iii) we use (ii):

$$
\begin{aligned}
& w(\nu+1)(*, \lambda q)=\bar{b}\left(\lambda q_{v}\right) \cdot \ldots \cdot \bar{b}\left(\lambda q_{1}\right)= \\
& \quad=\theta \bar{b}\left(q_{\nu}\right) \cdot \ldots \cdot \theta \bar{b}\left(q_{1}\right)=\theta\left(\bar{b}\left(q_{1}\right) \cdot \ldots \cdot \bar{b}\left(q_{\nu}\right)\right)=\theta w(\nu+1)(*, q) .
\end{aligned}
$$

4.2 Proposition. Let $\tilde{E}: P S^{2} \rightarrow \boldsymbol{R}$ be the energy on $P S^{2}$ where $S^{2}$ has the canonical metric. Let $\tilde{w}$ be any cycle homological to $w(v)$. Then

$$
\sup \tilde{E}|\tilde{w} \geqslant \sup \tilde{E}| w(v)=2 \pi^{2}(2 v+1)^{2}
$$

Proof. Among the curves of maximal $\widetilde{E}$-value in image $w(v)$ there is only one closed geodesic, i.e., the $(\nu+1)$-fold covering of the great circle $S_{0}^{1} \subset S^{2}$. It is the image of $\left(p_{0}, q_{0}\right)$ where $p_{0}=*^{\prime}=(-1,0,0)$ and $q_{0}$ is the $\nu$-tupel $(0, i)^{\nu}$. An application of the deformation $\tilde{D}(2$, ) will render the $\tilde{E}$-value of all other curves in image $w(\nu)<2 \pi^{2}(\nu+1)^{2}=\tilde{E}($ the $(\nu+1)$-fold covering of $S_{0}^{1}$ ).
$w(0)$ is a generator of $\pi_{1} \Omega S^{2}$. Thus, it cannot be homotopic to a $\tilde{w}$ having it image in the domain $\left\{\tilde{E}<2 \pi^{2}\right\}$ because then it would be homotopic to a map $S^{1} \rightarrow S^{2}$ and thus, null-homotopic.

Therefore it only remains to prove (4.2) for $v>0$. If there were a $\tilde{w} \sim w(v)$ with image $\tilde{w} \subset\left\{\tilde{E}<2 \pi^{2}(\nu+1)^{2}\right\}$ it would also be homotopic to a cycle representable as the negative bundle over some cycle $\beta$ of the space $B_{\mu+1} S^{2}$ of $(\mu+1)$-fold covered great circles, $\mu<\nu$. The dimension of such a cycle is the dimension of $\beta$ plus $2 \mu+1$. But as observed already before, the only non-zero cycles in $B_{\mu+1} S^{2}$ occur in dimensions 0 and 3 . $\operatorname{dim} \beta=0$ is impossible, since $\mu<\nu$. And $\operatorname{dim} \beta=3$ is also impossible since $(2 \mu+1)+$ $+3 \neq 2 v+1$.

So far, the product of circles $w(v)(p, q)=\bar{a}(b) \cdot \bar{b}\left(q_{1}\right) \cdot \ldots \cdot \bar{b}\left(q_{v}\right)$ is not necessarily in $P_{\nu} S^{2}$. We show that we can replace $w(v)$ by a homotopic element having this property:
4.3 Proposition. The cycle $w(v)$ is homotopic (modulo cells of codimension $\geqslant 2$ ) to a cycle belonging to $P_{\nu} S^{2} \cap\left\{\tilde{E} \leqslant 2 \pi^{2}(\nu+1)^{2}\right\}$. The new cycle which we denote again by $w(v)$ still has the property (4.1) (iii).

Proof. As we did observe already in the proof of (4.2), the cycle $w(\nu)$ is homologous to the non-trivial cycle given by the strong unstable manifold $W_{u u}\left(c^{\nu+1}\right)$ of a $(\nu+1)$-fold covered great circle $c^{\nu+1}$. This latter belongs to $P_{\nu} S^{2} \cap\left\{\tilde{E} \leqslant 2 \pi^{2}(\nu+1)^{2}\right\}$. To see this note that an eigenvector with the negative eigenvalue

$$
\lambda=4 \pi^{2}\left(\mu^{2}-(v+1)^{2}\right) /\left(1+4 \pi^{2} \mu^{2}\right)
$$

$\mu=0, \ldots, \nu$ is given by

$$
\begin{array}{ll}
X(t)=A \quad \text { with } \quad A \perp \dot{c}(t) & \text { for } \mu=0 \\
X(t)=A \cos 2 \pi \mu(t+\beta), & \text { for } 0<\mu \leqslant \nu+1, \beta \text { real. }
\end{array}
$$

$X\left(t_{0}\right)=X\left(t_{0}+x /(v+1)\right), x$ an integer, occurs at most $\nu$ times for a nonconstant $X(t)$. Thus, the curves near the origin of $W_{u u}\left(c^{\nu+1}\right)$ belong to $P_{\nu} S^{2}$. Along the $-\operatorname{grad} E$ flow lines the selfintersection number possibly increases. But that part of $\partial W_{u u}\left(c^{\nu+1}\right)$ which goes outside $P_{\nu} S^{2}$ is homologically negligable: Indeed, from (3.3) we have that the deformation $\widetilde{D}$ can be applied without encountering an obstruction since there are no critical points of selfintersection number $\nu$ below $\tilde{E}=2 \pi^{2}(\nu+1)$.

Consider now the cycle $w(v)$. The product $w(v)(q, p)$ of $v+1$ circles in general will have a selfintersection number $>\nu$. But if we apply the - grad $E$ deformation we will push this cycle into the Morse complex formed by the unstable manifolds of the $\mu$-fold covered great circles, $\mu \leqslant \nu+1$. Here, up to homologically negligable parts, $w(\nu)$ will be accomodated in $P_{\nu} S^{2}$, as we just saw. Since the $-\operatorname{grad} E$ deformation commutes with $\theta$, the relation (4.1) (iii) is being preserved.
5. - The cycles $w(\nu)$ constructed in paragraph (4) will yield a closed geodesic $c(v)$ in the usual manner, using the minimum-maximum method, cf. (5.5) below. We even will find a $c(v)$ with selfintersection number $v$. However, for the proof of the existence of infinitely many prime closed geodesics we need a second closed geodesic $c^{\prime}(\nu)$ with $v\left(c^{\prime}(\nu)\right)=\nu$.

This leads us to consider the mapping

$$
\begin{gathered}
w^{\prime}(\nu): D^{\prime}(\nu)=\left(S_{0}^{2}\right)^{\nu+1} \rightarrow \Omega S^{2} \\
q=\left(q_{0}, \ldots, q_{\nu}\right) \mapsto \bar{b}(q)=\bar{b}\left(q_{0}\right) \cdot \ldots \cdot \bar{b}\left(q_{\nu}\right)
\end{gathered}
$$

Note: $w^{\prime}(\nu)=w(v+1) \mid\{*\} \times\left(S_{0}^{2}\right)^{v+1}$. Moreover, by taking the modifications described in (4.3): image $w^{\prime}(v) \subset P_{\nu} M$.

We define the class $W^{\prime}(\nu)$ of mappings

$$
w^{\prime}: D^{\prime}(\nu) \rightarrow P M
$$

by the following conditions:
(i) image $w^{\prime} \subset P_{\nu} M$;
(ii) $w^{\prime}$ is homotopic to $w^{\prime}(\nu)$ by a sequence of admissible homotopies;
(iii) $w^{\prime}(\lambda q)=\theta w^{\prime}(q)$.

Here we mean by an admissible homotopy a continuous mapping

$$
h:[0,1] \times D^{\prime}(\nu) \rightarrow P_{\nu} M
$$

such that, if we put $h \mid\{\sigma\} \times D^{\prime}(\nu)=w^{\prime \sigma}, w^{\prime 0} \in W^{\prime}(\nu)$ and

$$
w^{\prime \sigma}(\lambda q)=\theta w^{\prime \sigma}(q)
$$

5.1 Proposition. Let $w^{\prime} \in W^{\prime}(v)$. Choose a $\tau_{0}$. Then

$$
w^{\prime \sigma}(q)=w^{\prime \sigma}\left(q_{0}, \ldots, q_{\nu}\right)=w^{\prime}\left(\gamma_{\tau_{0}+\pi \sigma} q_{0}, \ldots, \gamma_{\tau_{0}+\pi \sigma} q_{v}\right)
$$

is an admissible homotopy of $w^{\prime}$.
Proof. Immediate from $\lambda \gamma_{\tau}=\gamma_{\tau} \lambda$.
5.2 Proposition. Let $w^{\prime} \in W^{\prime}(\nu)$. Then there exists a $\widetilde{w}^{\prime} \in W^{\prime}(\nu)$, homotopic to $w^{\prime}$, such that

$$
E\left(\widetilde{w}^{\prime}(q)\right) \leqslant E\left(\widetilde{D}\left(2, w^{\prime}(q)\right)\right)
$$

for all $q \in D^{\prime}(\nu)$. Here $\widetilde{D}(2$,$) is a deformation from P_{\nu}^{*} M$ into $P_{\nu}^{\alpha} M$ of the type considered in (3). $x>0$ has been chosen so large that image $w^{\prime} \subset\{\boldsymbol{E} \leqslant \chi\}$.

Proof. Let $H_{0}^{2} \subset S_{0}^{2}$ be the half sphere given by $\left\{\left(\cos \alpha, e^{-i \varphi} \sin \alpha\right)\right.$; $0 \leqslant \alpha \leqslant \pi ; \pi / 2 \leqslant \psi \leqslant 3 \pi / 2\}$. Let

$$
f:\left(H_{0}^{2}\right)^{v+1} \rightarrow[0,1] \subset \boldsymbol{R}
$$

be a differentiable function with $f(q)=0$ for $q \in \partial\left(H_{0}^{2}\right)^{\nu+1}$ and $f(q)=1$ if every component of $q$ has distance $\geqslant \varepsilon$ from $\partial H_{0}^{2}$, some small $\varepsilon>0$.

Let $J(\nu+1)$ be the set of the $2^{\nu+1}(\nu+1)$-tupels $j=\left(j_{0}, \ldots, j_{\nu}\right) \in \mathbb{Z}_{2}^{\nu+1}$. If $q=\left(q_{0}, \ldots, q_{v}\right) \in\left(S_{0}^{2}\right)^{\nu+1}$ we define

$$
\lambda(j) q=\left(\lambda^{j_{0}} q_{0}, \ldots, \lambda^{j_{v}} q_{v}\right)
$$

Choose a fixed subset $J^{\prime}(\nu+1)$ of $J(\nu+1)$ of $2^{\nu}(\nu+1)$-tupels $j$. Then, for every $j \in J(v+1)$, either $j \in J^{\prime}(\nu+1)$ or else, $j+1=\left(j_{0}+1, \ldots, j_{v}+1\right) \in$ $\in J^{\prime}(\nu+1)$.

With this we define a homotopy $w^{\prime \sigma}$ of $w^{\prime}$ as follows, with $q^{*} \in\left(H_{0}^{2}\right)^{\nu+1}$ :

$$
w^{\prime \sigma}\left(\lambda(j) q^{*}\right)= \begin{cases}\tilde{\mathscr{D}}\left(2 \sigma f\left(q^{*}\right), w^{\prime}\left(\lambda(j) q^{*}\right)\right), & \text { if } j \in J^{\prime}(\nu+1) \\ \theta \tilde{\mathscr{D}}\left(2 \sigma f\left(q^{*}\right), w^{\prime}\left(\lambda(j+1) q^{*}\right)\right), & \text { if } j \notin J^{\prime}(\nu+1)\end{cases}
$$

This obviously is an admissible homotopy. Moreover, the inequality (5.2) is satisfied whenever $q \in D^{\prime}(v)$ can be written in the form $q=\lambda(j) q^{*}$ where $q^{*} \in\left(H_{0}^{2}\right)^{\nu+1}$ and $f\left(q^{*}\right)=1$.

To obtain (5.2) also for those $q=\lambda(j) q^{*}$ where $f\left(q^{*}\right)<1$ we apply to $\left(H_{0}^{2}\right)^{\nu+1}$ first one of the $(\nu+2)$ rotations

$$
\gamma_{\tau(\mu)}, \quad \tau(\mu)=\pi \mu /(\nu+2) ; \quad \mu=0, \ldots, \nu+1
$$

and follow this up with the previously defined homotopy. For every $q^{* *} \in\left(H_{0}^{2}\right)^{\nu+1}$ it will happen at least once that, if $\gamma_{\tau(\mu)} q^{* *}=\lambda(j) q^{*}, q^{*} \in$ $\in\left(H_{0}^{2}\right)^{v+1}$, then $f\left(q^{*}\right)=1$. This completes the proof of (5.2).

From our description of the Hopf map $b: S^{3} \rightarrow S^{2}$ in (4) we can see: If $H_{0}^{1} \subset S_{0}^{2}$ denotes the half great circle

$$
H_{0}^{1}=\{(\cos \alpha,-i \sin \alpha) ; 0 \leqslant \alpha \leqslant \pi\}
$$

then $b \mid H_{0}^{1}: H_{0}^{1} \rightarrow S^{2}$ has as image the great circle $S_{0}^{1} \subset S^{2}$. Actually, $b \mid H_{0}^{1}$ is injective, with the exception of the boundary points $* \in S^{3}$ and $*^{\prime} \in S^{3}$ of $H_{0}^{1}$ which both are mapped into $* \in \mathcal{S}_{0}^{1} \subset \mathbb{S}^{2}$.

Therefore we may write the cycle $\bar{a}: S_{0}^{1} \rightarrow \Omega S^{2}$ also in the form $\bar{b} \mid H_{0}^{1}: H_{0}^{1} \rightarrow \Omega S^{2}$. Consequently we denote also $H_{0}^{1} \times\left(S_{0}^{2}\right)^{\nu}$ by $D(\nu)$. Then

$$
w(v)=w^{\prime}(v) \mid D(v)
$$

This leads us to define the class $W(\nu)$ as being mappings

$$
w: D(v) \rightarrow P M
$$

where $w=w^{\prime} \mid D(\nu), w^{\prime} \in W^{\prime}(\nu)$. Clearly, any $w \in W(\nu)$ is homotopic to $w(\nu)$.

Let $w \in W(v)$. We call $(p, q) \in D(v)$ a regular point if it is in the domain of a non-degenerate singular simplex of the map $w: D(\nu) \rightarrow P M$, considered as sum of singular simplices. Similarly we define a regular point of $w^{\prime} \in W^{\prime}(\boldsymbol{v})$.
5.3 Proposition. The set of regular points of $w \in W(\nu)$ belonging to $w^{-1}\left(P_{\nu} M-P_{\nu-1} M\right)$ is non-empty.

Proof. We can assume: $v>0$. Since $w$ is a non-trivial cycle there are regular points. If $w^{-1}\left(P_{\boldsymbol{v}} M-P_{\nu-1} M\right)$ were empty, i.e., if image $w \subset P_{\nu-1} M=$ $=P_{\nu-1} S^{2}$, then image $w$ could be deformed in $P_{\nu-1} S^{2}$ below the $\widetilde{E}$-level $2 \pi^{2}(2 v+1)^{2}$, since there are no closed geodesics in $P_{\nu-1} S^{2}$ with $\tilde{E}$-value $>$ $>2 \pi^{2}(2 v-1)^{2}$. But this contradicts (4.2). The degenerate simplices in $w$ are negligable homologically. Therefore, since $w$ is non-trivial, the same argument shows that there are non-degenerate simplices in $P_{\nu} M-P_{\nu-1} M$.
5.4 Proposition. The set of regular points of $w^{\prime} \in W^{\prime}(\nu)$ having image in $\boldsymbol{P}_{\boldsymbol{\nu}} \boldsymbol{M}-\boldsymbol{P}_{\boldsymbol{\nu}_{-1}} M$ is non-empty.

Proof. Since $w^{\prime} \in W^{\prime}(\nu)$ is not a non-trivial rational cycle we cannot employ the same arguments as for the proof of (5.3). Instead, we will use the fact that a $w^{\prime} \in W^{\prime}(\nu)$ determines a non-trivial $\boldsymbol{Z}_{2}$-cycle of the space $P M / \theta$ of unoriented parameterized closed curves on $M$. Here, $P M / \theta$ is the quotient of $P M$ with respect to the $Z_{2}$-action generated by the orientation reversing map $\theta$. Note: $\theta P_{\nu} M=P_{\nu} M$.

To make this more precise we consider the quotient space $D^{\prime}(\nu) / \lambda$ of $D^{\prime}(\nu)$ with respect to the $\boldsymbol{Z}_{2}$-action generated by $\lambda: D^{\prime}(\nu) \rightarrow D^{\prime}(\nu)$, cf. (4). Using the notations employed in the proof of (5.2) we can represent the fundamental $\boldsymbol{Z}_{2}$-homology class of the space $D^{\prime}(\nu) / \lambda$ by the cycle

$$
z_{2 v+2}: Z^{2 v+2}=\bigcup_{j \in J^{\prime}(v+1)} \lambda(j)\left(H_{0}^{2}\right)^{v+1} \rightarrow \overrightarrow{l \lambda} D^{\prime}(v) / \lambda
$$

The 1-dimensional $\boldsymbol{Z}_{2}$-homology class is being represented by the cycle

$$
z_{1}: Z^{1}=H_{0}^{1} \rightarrow D^{\prime}(\nu) / \lambda
$$

Here we have considered $H_{0}^{1}$ as subset of the first factor of $D^{\prime}(v)=\left(S_{0}^{2}\right)^{\nu+1}$. Finally, the 1-codimensional $\boldsymbol{Z}_{2}$-homology class of $D^{\prime}(\boldsymbol{v}) / \lambda$ is represented by the cycle

$$
z_{2 v+1}: Z_{2 \nu+1}=H_{0}^{1} \times Z^{2 v} \rightarrow \vec{\nmid} D^{\prime}(\nu) / \lambda
$$

with $H_{0}^{1} \times Z^{2 v} \subset S_{0}^{2} \times\left(S_{0}^{2}\right)^{\nu}=D^{\prime}(\nu)$.
$D^{\prime}(v) / \lambda$ has the important property that any 1-cycle and any $(2 v+1)$ cycle, when in general position, intersect in an odd number of points. That is to say, $z_{2 \nu+1}$ represents the 1-dimensional $\boldsymbol{Z}_{2}$-cohomology class of $D^{\prime}(\nu) / \lambda$. Note that $\left(z_{1}, z_{2 v+1}\right)$ are not in general position.

An element $w^{\prime} \in W^{\prime}(\nu)$ can be viewed as the 2 -fold covering of a mapping

$$
w^{\prime} / \lambda: D^{\prime}(v) / \lambda \rightarrow P_{v} M / \theta
$$

Assume now that $w^{\prime}$ has no regular points. Since $w^{\prime} \mid D(v) \in W(v)$ is a non-trivial cycle this implies that some 1-cycle $z_{1}^{\prime} \sim z_{1}$ becomes a trivial cycle under $w^{\prime} \mid z_{1}^{\prime}$. But this is impossible since $w^{\prime} \mid z_{1}^{\prime}$ is homologous to $w(0)$ which is a non-trivial cycle, representing the generator of $H_{1} \Omega S^{2}=H_{1} P S^{2}$, cf. (3).

That the image of the interior of the set of regular points meets $P_{\nu} M-P_{\nu-1} M$ is proved now just as in (5.3).

We now can define:

$$
\begin{aligned}
& \varkappa(v)=\inf _{w \in W(v)} \sup \{E(c) ; v(c)=v, c=w(p, q),(p, q) \text { regular point }\} \\
& \varkappa^{\prime}(v)={ }_{w^{\prime} \in W^{\prime}(v)} \sup \left\{E(c) ; v(c)=v, c=w^{\prime}(q), q \text { regular point }\right\}
\end{aligned}
$$

5.5 Theorem. For every $\nu=0,1, \ldots$, there exist closed geodesics $\left\{c(\nu), c^{\prime}(\nu)\right\}$ with selfintersection number $\nu$ and $E(c(v))=x(\nu), E\left(c^{\prime}(\nu)\right)=\varkappa^{\prime}(\nu)$. Moreover

$$
\begin{equation*}
0<\varkappa(v) \leqslant \varkappa^{\prime}(v)<\varkappa(v+1) . \tag{*}
\end{equation*}
$$

Here equality can hold only, if there are infinitely unparameterized closed geodesics with $E$-value $\varkappa(\nu)=\varkappa^{\prime}(\nu)$ and selfintersection number $\nu$.

Proof. Since $w \in W(v)$ contains as restriction $w \mid H_{0}^{1} \times\{* \ldots *\}$ an element of $W(0), x(v)>0$ will follow if we show: $x(0)>0$. But this can easily be seen as follows: If $\varkappa(0)=0$ then there would exist a $w \in W(0)$ with image $w \subset\left\{E \leqslant \eta^{2} / 2\right\}, \eta>0$ as in (3). For every $p \in H_{0}^{1}$, the closed curve $w(p)$ therefore can be retracted along the radii of the convex $2 \eta$-neighborhood around $w(p)(0)=w(p)(1)$ into $w(p)(0)$. Thus, $w(0)$, being homotopic to $w$, becomes homotopic to a map $S^{1} \rightarrow M=S^{2}$, i.e., to a constant map, which is a contradiction.

Let $q \in D^{\prime}(\nu)$ be a regular point of a $w^{\prime} \in W^{\prime}(\nu)$ : Then there exists a $w \in W(v)$ containing $q$ among its domain of definition. Thus, $q$ is regular point also for some $w \in W(v)$. This shows: $x(v) \leqslant \varkappa^{\prime}(\nu)$.

To prove the existence of a closed geodesic $c(v)$ with $E(c(v))=\varkappa(\nu)$,
$\nu(c(v))=\nu$ we observe that we get from the definition of $\varkappa(\nu)$ a sequence $\left\{w_{n}\right\}$ in $W(\nu)$ and a sequence $\left\{r_{n}\right\}$ in $D(\nu), r_{n}$ interior regular point of $w_{n}$, such that for $c_{n}=w_{n}\left(r_{n}\right): \nu\left(c_{n}\right)=\nu$ and

$$
x(\nu) \leqslant E\left(\widetilde{D}\left(2, c_{n}\right)\right) \leqslant E\left(c_{n}\right)
$$

and $\lim E\left(c_{n}\right)=\varkappa(v)$. According to (3.3) we can assume that $\lim c_{n}$ exists and is a closed geodesic $c(v)$. We can assume that also $\mathfrak{D}\left(2, c_{n}\right)$ is in the image of an interior regular point of $\widetilde{\mathscr{D}}\left(2, w_{n}\right) \in W(\nu)$ and belongs to $P_{\nu} M-P_{\nu-1} M$. (3.3) then implies: $v(c(\nu))=\nu$.

The same arguments lead to the existence of a $c^{\prime}(\nu)$ with the desired properties.

To prove $\varkappa^{\prime}(\nu)<\varkappa(\nu+1)$ we first show that $\varkappa^{\prime}(\nu) \leqslant \varkappa(\nu+1)$. To see this we deduce from our previous arguments the existence of a sequence $\left\{w_{n}\right\} \in$ $\in W(\nu+1)$, hence, $\left\{w_{n}^{\prime}:=w_{n}\left\{\{*\} \times D^{\prime}(\nu)\right\} \in W^{\prime}(\nu)\right.$, and of a $q \in D^{\prime}(\nu)$ being regular for all $w_{n}^{\prime}$ such that the sequence $\left\{c_{n}=w_{n}^{\prime}(q)=w_{n}(*, q)\right\}$ converges to $c^{\prime}(\nu)$ with $\nu\left(c_{n}\right)=\nu$.

We claim that for any $q \in D^{\prime}(\nu)$, regular for $w_{n}^{\prime}$, there exist $p \in H_{0}^{1}$, arbitrarily close to $* \in H_{0}^{1}$ such that $(p, q) \in D(\nu+1)$ is regular for all $w_{n}$. This means that $(*, q) \in D(v+1)$ is a boundary regular point for $w_{n}$. Indeed, the cycle $w_{n} \mid H_{0}^{1} \times\{q\}$ is homotopic to $w_{n} \mid H_{0}^{1} \times\{*\} \in W(0)$ and hence nontrivial. Therefore, after possibly some homotopic modification of $w_{n}$ near $(*, q), w_{n} \mid H_{0}^{1} \times\{q\}$ will be regular at $(*, q)$, thus our claim.

Consider now $w_{n} \mid\{p\} \times D^{\prime}(\nu)$. As long as the image is in $P_{\nu} M$, while $p \in \boldsymbol{H}_{0}^{1}$ moves away from $*$, we get an element of $W^{\prime}(\nu)$ and thus, $\sup E \mid\left(w_{n} \mid\{p\} \times\right.$ $\left.\times D^{\prime}(\nu)\right) \geqslant \varkappa^{\prime}(\nu)$. But it is impossible that image $w_{n} \mid\{p\} \times D^{\prime}(\nu) \subset P_{\nu} M$, for all $p \in H_{0}^{1}$, cf. (5.3). Hence, there is one $p_{0} \in H_{0}^{1}$ (possibly $p_{0}=*$ ) such that $p_{0}$ is the limit of points $p \in H_{0}^{1}$ with image $w_{n} \mid\{p\} \times D^{\prime}(\nu) \cap P_{v+1} M \neq \emptyset$. Thus, $x(\nu+1) \geqslant x^{\prime}(v)$.

To see that actually $\varkappa^{\prime}(\nu)<\varkappa(v+1)$ we note that the set $C(v)$ of closed geodesics $c$ with $E(c)=\varkappa^{\prime}(\nu)$ and $\nu(c)=\nu$ and the set $C(\nu+1)$ of closed geodesics $c$ with $E(c)=\varkappa^{\prime}(\nu)$ and $\nu(c)=\nu+1$, if not empty, are disjoint compact sets. Since we just showed that in every neighborhood of $c(\nu) \in$ $\in C(\nu)$ there are elements $c_{n}=w_{n}(p, q) \in P_{\nu+1} M-P_{\nu} M,(p, q)$ regular for $w_{n} \in W(v+1)$, we see that $\varkappa(v)=\varkappa(v+1)$ is impossible.

It remains to discuss the case $\varkappa(\nu)=\varkappa^{\prime}(\nu)=$ (briefly) $\varkappa_{0}$. We want to derive a contradiction from the assumption that the set $C$ of closed geodesics $c$ with $E(c)=\varkappa_{0}, \nu(c)=\nu$ consists of only finitely many $\mathcal{S}$-orbits $\mathcal{S} . c$. In this case we could choose arbitrarily small open neighborhoods $\mathcal{U}$ of $C$ of the following type: $\mathcal{U}$ is the union of finitely many pairwise disjoint $S$-invariant open neighborhoods $\mathcal{U}(S . c)$ of the finitely many $\mathbb{S}$-orbits $\mathbb{S} . \boldsymbol{c}$. We
also may assume:

$$
\text { U(S. } \theta c)=\theta ひ(S . c)
$$

From (3.4) and (5.2) we have the existence of an $\varepsilon>0$ and a $w^{\prime} \in W^{\prime}(v)$ such that

$$
\text { image } w^{\prime} \subset \text { U } \cup P_{\nu}^{x_{0}-\varepsilon} M
$$

$x(v)=x^{\prime}(v)$ means that the image under $w^{\prime}$ of every cycle $z_{2 v+1}^{\prime} \sim z_{2 v+1}$ (cf. the proof of (5.4)) meets $\mathcal{U} / \theta$. Or equivalently: The carrier of every non-trivial 1-dimensional $\boldsymbol{Z}_{2}$-cocycle meets $\mathcal{U} / \theta$. That is to say, there is a $z_{1}^{\prime} \sim z_{1}$ where $w^{\prime}\left(z_{1}^{\prime}\right) \subset \mathcal{U}(S . c) / \theta$, some $S$-orbit $S . c$ in $C$. But $S . c$ can be retracted into a trivial $\mathcal{S}$-orbit $\subset P^{0} M=M$, thus, the cycle $w^{\prime} \mid z_{1}^{\prime}$ is trivial which is a contradiction since it is homologous to the cycle $w(0)$.

This completes the proof of (5.5).
6. - We can now prove our main result.
6.1 Theorem. Let $M$ be a closed surface of genus 0 . Then there exist on $M$ infinitely many unparameterized prime closed geodesics.

Notes. 1) Lusternik and Schnirelmann 1.c. had proved the existence of three closed geodesics without selfintersections on an orientable surface of genus 0 . The ellipsoid with three pairwise different axes, all approximately of the same length, gives an example of such a surface where there exist exactly three closed geodesics without selfintersections, i.e., the three principal ellipses of the ellipsoid.

As was shown by Morse [Mo], cf. also LCG, the next prime closed geodesic in $E$-value and hence in length, after these three relatively short closed geodesics, can have an $E$-value greater than any prescribed number, if only the three axes have their length sufficiently near 1. This fourth prime closed geodesic therefore also will have arbitrarily large selfintersection number.
2) The theorem is a special case of our theorem (4.3.5) in LCG which states that on every compact riemannian manifold with finite fundamental group there exist infinitely many prime closed geodesics. The proof of this more general theorem uses various deep results, among others the GromollMeyer theorem [GM], Sullivan's theory of the minimal model [ Su ] and the structure of the Morse complex of the Hilbert manifold of closed $H^{1}$-curves with respect to a non-degenerate energy function $E$, cf. LCG for details.

In contrast, the present proof of this result for a closed surface of genus 0 can be considered elementary. We could have shortened it were it not for an attempt to keep the topological prerequisites to a minimum.

Proof. We can assume $M$ is orientable. That is $M$ is the 2 -sphere with an arbitrary riemannian metric. Indeed, otherwise the universal covering $\tilde{M}$ of $M$ is of this type. If we have infinitely many prime closed geodesics on $\tilde{M}$ then their projection into $M$ under the covering map $\tilde{M} \rightarrow M$ will yield infinitely many prime closed geodesics on $M$ also since the image of different prime closed geodesics has different underlying prime closed geodesics.

We now will derive a contradiction from the assumption that there are on $M$ only finitely many $S$-orbits of prime closed geodesics. Let $\left\{c_{1}, \ldots, c_{s}\right\}$ representatives of these different $S$-orbits.

From (5.5) we get for every $v=0,1, \ldots$ a pair $\left\{c(\nu), c^{\prime}(\nu)\right\}$ of closed geodesics with selfintersection number $\nu$. Under our assumption the relations (*) in (5.5) become

$$
0<E(c(v))<E\left(c^{\prime}(\nu)\right)<E(c(\nu+1)) .
$$

For a fixed $\nu$ there exist integers $a(\nu), b(\nu), 0 \leqslant a(\nu)<b(\nu) \leqslant s$, such that both, $c(\nu+a(\nu))$ and $c(\nu+b(\nu))$ have the same underlying prime closed geodesic. Here and in the following we identify geodesics which lie on the same $S$-orbit.

Taking iterated subsequences we see: There exists a strictly increasing sequence $\{v(k)\}$ of integers $\geqslant 0$ and $a, b, 0<b \leqslant s$, such that, for all $k \in N$, $c(v(k))$ and $c(v(k)+b)$ have the same underlying prime closed geodesic, say $c$. Moreover, also the geodesics $c^{\prime}(v(k))$ all have the same underlying prime closed geodesic, say $\boldsymbol{c}^{\prime}$.

Put $E(c)=\varkappa \quad E\left(c^{\prime}\right)=\varkappa^{\prime}$. Denote the multiplicity of $c(\nu(k)), c^{\prime}(\nu(k))$ and $c(\nu(k)+b)$ by $m(k), m^{\prime}(k)$, and $\tilde{m}(k)$, respectively. Put $\nu(c)=\nu$, $\nu\left(c^{\prime}\right)=\nu^{\prime}$. Then we have the following relations, cf. (5.5), (2.1):

$$
\begin{equation*}
m(k)^{2} \varkappa<m^{\prime}(k)^{2} \varkappa^{\prime}<\tilde{m}(k)^{2} \varkappa \tag{*}
\end{equation*}
$$

$$
\left\{\begin{array}{l}
\text { (i) } v(k)+1=m(k)^{2} v+m(k)=m^{\prime}(k)^{2} v^{\prime}+m^{\prime}(k)  \tag{**}\\
\text { (ii) } v(k)+b+1=\tilde{m}(k)^{2} v+\tilde{m}(k) .
\end{array}\right.
$$

This shows that with $k \rightarrow \infty$ also $m(k), m^{\prime}(k), \tilde{m}(k) \rightarrow \infty$. Moreover:
$(* * *) \quad \lim _{k} m(k)^{2} / \tilde{m}(k)^{2}=1 ; \quad \lim _{k} m^{\prime}(k)^{2} / m(k)^{2}=x / \varkappa^{\prime}$.
This follows from ( $*$ ) together with $(* *)$.
If $v=\nu^{\prime}=0, m(k)=m^{\prime}(k), \varkappa=\varkappa^{\prime}$, which contradicts $(*)$.
If $\nu>0$ we have from

$$
1=\frac{m(k)^{2}(v+1 / m(k))}{m^{\prime}(k)^{2}\left(v+1 / m^{\prime}(k)\right)}
$$

and ( $* * *$ ):
$(* * * *) \quad x / x^{\prime}=\nu / \nu^{\prime}$.
From (*) and (**) follows:

$$
\begin{aligned}
& m(k)^{2} v<m^{\prime}(k)^{2} v^{\prime}<\tilde{m}(k)^{2} v \\
& \tilde{m}(k)-b<m^{\prime}(k)<m(k)
\end{aligned}
$$

Thus, with $(* * *), x=x^{\prime}, \nu=\nu^{\prime}$, which again contradicts ( $*$ ). This completes the proof.

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