Physics. - "On the determination of quanta-conditions by means of adiabatic invariants." By G. Krutkow. (Communicated by Prof. P. Ehrenfest.)
(Communicated at the meeting of September 25, 1920).
In a series of papers Ehrenfest has shown that only such functions of the general co-ordinates of a mechanical system may be quanticized as are adiabatic invariants ${ }^{1}$ ). These functions can always be found ${ }^{2}$ ). Moreover, as we shall see, theory may answer the question as to the number of essential adiabatic invariants, which in accordance with the quanta-hypothesis have to assume discontinuous values. If we suppose that the "density of probability" of the motion of the system, when not adiabatically acted upon, does not depend explicitly on the time, and if then by means of some hypothesis or some theorem which is derived from the properties of the system, we replace the time-mean of a phase-function by a numerical mean, it follows immediately that the number of essential invariants is equal to the number of determining quantities of the system which is left after the numerical mean has been determined (comp. equations (12) sqq. below). By the determination of the adiabatic invariants and the separation of the essential ones the uncertainty as to the choice of the forms of motion which are admissible on the quanta-hypothesis, becomes materially lessened. Still we must not expect that the adiabatic invariants which we have found are necessarily those which have to be quanticized: any arbitrary function of those quantities is again an adiabatic invariant and has thus equal claims to being selected. However, this liberty of choice can be somewhat restricted; there is a further condition to which we may subject the quanta-functions. This condition is of the nature of a hypothesis, but we may give it a simple statistical interpretation. In every case, where the theory of quanta has been applied with success ${ }^{2}$ ), the condition is fulfilled. It was introduced by Planck as a fundamental theorem for a complete determination of the quantities
${ }^{1}$ ) P. Ehrenfest. These Proc. XIX N0: 3, p. 576. Ann. der Phys. 51 (1916) p. 327.
${ }^{2}$ ) G. Krutrow. Proc. Amst XXI p. 1112. 1919.
3) My knowledge of the literature of the subject does not, however, extend beyond the beginning of 1917.
which have to be quanticized ${ }^{1}$ ). A new proof will be given by establishing a connection between the adiabatic invariants and the phase-space (below 18').
This connection, which will be found to arise in a natural way, with a concept derived from statistical mechanics, strengthens the bond between it and the theory of quanta, a bond which, as it seems to me, has gone into the background in the latest development of the theory or at least has not been sufficiently emphasized, although in my opinion it is of great importance. In view of this connection I think that the only justification of the expression "action-quantum" is the fact that it recalls to our mind the dimensions of the phase-extension
Another conception of great importance to the theory of quanta which will find a place in our classification is Planck's ${ }^{3}$ ) coherence of degrees of freedom. To me it seems of fundamental importance. Its meaning will be found to appear very clearly by a juxtaposition of the properties of a conditionally periodic system and a Boltzmann "ergode".
This coherence of degrees of freedom must be very clearly distinguished from what is called "degeneration" ${ }^{3}$. For instance from our point of view an ergodic system is to the highest degree coherent, but could in no case be called degenerated. For a degenerated system the number of essential adiabatic invariants is greater than that of the degrees of freedom, for a coherent system it is smaller.

The question arises: must the supernumerary adiabatic invariants of a degenerated system be quanticised or, as suspected by Schwarzschid ${ }^{4}$ ), is the number of quanta-conditions smaller for such a system for the normal case without degeneration?

For the solution of these questions the three steps which have been taken viz. (1) establishment of the adiabatic invariants (2) selection of the essential ones and (3) "normalisation" of the latter are insufficient. In order to get nearer to the solution we must, I think, take into account, that the quanta-functions must liave a meaning which is independent of the system of co-ordinates. We may undoubtedly postulate this: if the quanta-laws are really physical laws, they must necessarily satisfy this condition. The question is, how to formulate this new invariance of the quanta-functions? I shall not try to discuss it here in general; but only remark that

[^0]${ }^{4}$ ) K. Schwarzschild, I.c
we may return from the canonical equations which are so convenient in the theory of quanta to the equations on cartesian co-ordinatos. Here the invariance in question means: invariance with respect to the groups of rotations and translations; vector-analysis thus provides the means of testing hypothetical quanta-quantities for the new postulata ${ }^{1}$ ).

The above mentioned means enable us in special cases to separate the quanta-quantities without ambiguity, for instance for the mechanical systems considered by Planck in the paper quoted. In some cases, however, an ambiguity remains, which we may get rid of in the following manner: by putting all but one of the quanta-quantities equal to nought, a "singular motion" must be obtained. In this manner we are able to make a connection between the methods sketched out above and Planck's theory on the physical structure of the phase-space, Planck's singular motions forming the last step in the series. We may recapitulate as follows:

The quanta-quantities are $(\alpha)$ functions of the integrals of the equations of motion ( $\boldsymbol{\beta})$ adiabatic invariants which $(\gamma)$ must be "normalized" and ( $\delta$ ) have a meaning which is independent of the system of co-ordinates and finally (e) yield singular motions in Planck's sense of the expression.

## §1. The fundamental equation.

Let a mechanical system of $n$ degrees of freedom be given by its canonical equations of motion

$$
\begin{equation*}
\dot{p}_{i}=-\frac{\partial H}{\partial q_{i}} ; \dot{q}_{i}=\frac{\partial H}{\partial p_{i}}(i=1,2, \ldots . n) \tag{1}
\end{equation*}
$$

We shall consider a number of systems and introduce a function $\varrho\left(p_{i}, q_{i}, t\right)$ which may be called the density of probability: $\rho$ must satisfy the fundamental equation of statistical mechanies ${ }^{2}$ ):

$$
\begin{equation*}
\frac{\partial \varrho}{\partial t}+\sum_{i=1}^{n}\left(\frac{\partial \varrho \dot{p}_{i}}{\partial p_{i}}+\frac{\partial \dot{\varphi}_{i}}{\partial q_{i}}\right)=0 \tag{2}
\end{equation*}
$$

or, using (1):

$$
\frac{\partial \varrho}{\partial t}+\sum_{i=1}^{n}\left(\frac{\partial \varrho}{\partial p_{i}} \dot{p}_{i}+\frac{\partial \varrho}{\partial q_{i}} \dot{q}_{i}\right)=\frac{d \varphi}{d t}=0
$$

0 is therefore a function of the integrals of the equation (1).

[^1]${ }^{2}$ ) J. W. Gibbs. Scientific papers. II' p. 16; Statistical Mechanics. Chapter I.

If we suppose that the condition is stationary

$$
\begin{equation*}
\frac{\partial \varrho}{\partial t}=0 \tag{3}
\end{equation*}
$$

it follows that: $\varrho$ is a function of those integrals which do not contain $t$ explicitly, i.e. of ( $2 n-1$ ) integrals, if only, as we shall suppose all the time, $H$ does not depend on $t$ explicitly.

We are at liberty to understand by $\rho$ the density of probability a posteriori or a priori. When applied to the theory of quanta our result expresses the fact, that the quanta-quantities are functions of the (2n-1) integrals of equations (1) which are independent of $t$.
Replacing the $2 n$-dimensional phase-space ( $p_{i}, q_{i}$ ) by the corresponding integral space $\left.\left(c_{i}, t_{i}\right)^{1}\right)$ the "path" of the system is a straight line parallel to the $t_{1}$-axis. We can describe these lines either by making $t$ increase, i.e. by following a definite system in its motion, or by keeping $t$ constant and varying $r$, i.e. considering together all the systems with given $c_{1} \ldots c_{n} t_{2} \ldots t_{n}$ and all possible values of $t_{1}$.

## §2. $H$ contains a variable parameter.

If $H$ contains a parameter which may either have a constant value as in the case just considered or vary slowly ${ }^{2}$ ), the quantities $c_{3}$ and $t_{1}$ are no longer constant, but variable; they have to satisfy the following "equations of motion" ${ }^{8}$ ):

$$
\begin{equation*}
\dot{c}_{i}=-\frac{\partial K}{\partial t_{i}} ; \quad \dot{t}=\frac{\partial K}{\partial c_{i}} \tag{3}
\end{equation*}
$$

where:

$$
\begin{equation*}
K=c_{1}+\left(\frac{\partial V}{\partial a} \dot{a}\right) \tag{4}
\end{equation*}
$$

${ }^{1}$ ) As in a previous communication we write the integrals of equations (1) in the form:

$$
H=c_{1}, \quad H_{2}=e_{2} ; \ldots, \quad H_{n}=e_{n}
$$

and

$$
\frac{\partial V}{\partial c_{1}}=t+t=t_{1}, \quad \frac{\partial V}{\partial c_{9}}=t_{2}, \ldots, \quad \frac{\partial V}{\partial c_{n}}=t_{n}
$$

where $c_{1}, \ldots, c_{n}, \tau, t_{2}, \ldots, t_{n}$ represent the $2 n$ integration-constants and $V$ JAcobl's characteristic function. Comp. Proc. XXI, p. 1112, 1919.
2) The slowness is expressed in the fact, that $H$ contains only $a$, not the correspondiug momentum.
${ }^{3}$ ) Proc. 1.c.
or putting $\dot{a}=$ const. approximately and representing the derivatives with respect to $a$ by dashes:

$$
\left.\begin{array}{c}
c_{i}^{\prime}=-\frac{\partial}{\partial t_{i}}\left(\frac{\partial V}{\partial a}\right) \quad ; \quad t^{\prime} x=\frac{\partial}{\partial c_{x}}\left(\frac{\partial V}{\partial a}\right)  \tag{}\\
t_{1}^{\prime}=\frac{1}{a}+\frac{\partial}{\partial c_{1}}\left(\frac{\partial V}{\partial a}\right)
\end{array}\right\}\binom{i=1,2, \ldots, n}{x=2,3, \ldots, n}
$$

Since the equations have the canonical form, we have as the fundamental equation :

$$
\begin{equation*}
\frac{\partial \varrho}{\partial a}+\sum_{i=1}^{n}\left(\frac{\partial \varrho}{\partial c_{i}} o_{i}^{\prime}+\frac{\partial \varrho}{\partial t_{i}} t_{i}^{\prime}\right)=0 . \tag{5}
\end{equation*}
$$

Here we may not as before take $\frac{\partial \varrho}{\partial \alpha}=0$. A further difficulty presents itself: starting from a special line parallel to the $t_{1}$-axis in the ( $c_{i}, t_{i}$ ) space - a special "stream-line" - if we now vary $a$, as equation (3) or ( $3^{\prime}$ ) show, the stream-line becomes broken up. If we then keep a constant again and take together the points, that lie on a straight line, $\varrho$ will vary along this stream-line, since it contains points of different origin. Thus on the new line $\rho$ is not stationary, but explicitly dependent on $t_{1}$

We now form ${ }^{1}$ ) the time-mean of $\rho$, which we shall call $\varrho$ and the difference $\varrho-\bar{\varrho}$. Since $\int d t_{1}(\rho-\bar{Q})=0$, the quantity $\varrho-\bar{\varrho}$ in its dependence on $t_{1}$ shows elevations and depressions round about $\bar{\rho}$, the sum of the surfaces of the former being equal to that of the latter. Each point carries its $a-\bar{\rho}$ value along with it and hence the curve shifts regularly with the time $t_{1}$. A stationary curve represents the tendency towards condensation (in an elevation) or rarefaction (in a depression) for the points of the stream-line, on the supposition of the change of $a$ being sufficiently small. If we make our moving curve slide along the stationary one, in the course of time elevations will cover depressions "and vice versa. A further small change of a may therefore produce a diminution of the difference $\rho-\bar{Q}$. By this reasoning it becomes clear that starting from a stationary density a sufficiently slow change of $a$ will to a corresponding degree of approximation produce a stationary density ${ }^{2}$ ).

[^2]We will therefore suppose that a changes slowly in the sense of the theory of adiabatic invariants. Let $D_{a}$ be the total change of $a$, i.e.

$$
D_{a}=\dot{a} \int d t
$$

and let $D c_{i}$ and $D t_{x}$ represent the corresponding changes of $c_{i}$ and $t_{x}$; considering further that

$$
\dot{a}=\mathrm{const}
$$

and hence

$$
\dot{a}=\frac{D_{a}}{D t} \quad \text { with } \quad\left(D t=\int d t\right):
$$

we find

$$
\begin{equation*}
\frac{D_{c_{i}}}{D a}=-\frac{\partial}{\partial t^{i}}\left(\frac{\partial V}{\partial a}\right) \quad ; \quad \frac{D t_{x}}{D a}=\frac{\partial}{\partial c_{x}}\left(\frac{\partial V}{\partial a}\right), \tag{6}
\end{equation*}
$$

where the horizontal line indicates the time-mean. If in equation (5) we take $c^{\prime}{ }_{i}$ and $t^{\prime}{ }_{i}$ to mean these time-means, we obtain

$$
\begin{equation*}
\frac{\partial \varrho}{\partial a}+\sum_{i=1}^{n}\left(\frac{\partial \varrho}{\partial c_{i}}{\frac{D c_{i}}{D a}}_{i}+\frac{\partial \varrho}{\partial t_{i}} \frac{D t_{i}}{D a}\right)=0 \tag{7}
\end{equation*}
$$

Since $\rho$ is independent of $t$, the corresponding term under the summation-sign in (7) mnst be omitted.
§3. Phase-space and adiabatic invariants.
The stationary density $\rho$ need not depend on all the variables

$$
c_{1}, \ldots, c_{n} ; t_{2}, \ldots ., t_{n}
$$

For example in a conditionally-periodic system without commen surable relations between the periodicity-moduli $\rho$ depends on the quantities $c_{i}$ only. This follows from the theorem which allows us to replace the time-average by an averaging over a $\Omega$-cell ${ }^{1}$ ). For an ergodic (or quasi ergodic) system in consequence of the ergodehypothesis $\rho$ depends on the energy $c_{1}$ only. We shall here suppose, that $\varrho$ depends on $k$ quantities $(k \leqq n)$, which we shall indicate by

$$
c_{1}, c_{3}, \ldots, c_{k}
$$

These integrals may be called essential integrals. Our supposition with regard to o comes to the same as assuming that for our system the time-mean may be replaced by a definite numerical mean. To compute this we proceed as follows.

Suppose the system of equations

$$
\begin{equation*}
H_{1}=c_{1}, H_{2}=c_{3}, \ldots, H_{k}=c / 6 \tag{8}
\end{equation*}
$$

to be soluble with respect to $p_{1}, p_{9}, \ldots, p_{k}$, thus
${ }^{1}$ ) J. M. Burgers. 1.c. and my paper 1.c.

$$
p_{\lambda}=k_{\lambda}\left(q_{2}, \ldots, q_{n}, p_{k+1} \cdot p_{n} ; c_{1}, \ldots c_{k}\right)(\lambda=1,2, \ldots, k)
$$

Introducing the differentials $d c_{1} \ldots . d c_{c}$ instead of the differentials $d p_{1} \ldots d p_{k}$ into the phase-integral

$$
\begin{equation*}
I=\iint \ldots \iint d p_{1} \ldots d p_{n} d q_{1} \ldots d q_{n} \tag{9}
\end{equation*}
$$

we find

$$
\begin{equation*}
I=\int \ldots \int d c_{1} \ldots d c_{k} \int \ldots \int d p_{k+1} \ldots d p_{n} d q_{1} \ldots d q_{n} \frac{\partial\left(p_{1}, \ldots p_{k}\right)}{\partial\left(c_{1}, \ldots c_{k}\right)} \tag{10}
\end{equation*}
$$

or

$$
I=\int \ldots \int d c_{1} \ldots d c_{k} \omega\left(c_{1}, \ldots c_{k}\right)
$$

where

$$
\begin{equation*}
\omega\left(c_{1}, \ldots, c_{k}\right)=\int \ldots \int d p_{k+1} \ldots d q_{n} \frac{\partial\left(p_{1}, \ldots, p_{k}\right)}{\partial\left(c_{1}, \ldots, c_{k}\right)} \tag{11}
\end{equation*}
$$

In (11) the integration has to be carried out, the limits being determined by (8). From the ( $p_{i}, q_{i}$ )-space or the ( $c_{i}, t_{i}$ )-space we may pass to the $k$-dimensional ( $c_{i}, \ldots c_{k}$ ) space. A streamline of the former space corresponds to a fixed point in the latter. The density $\rho$ is replaced by $\rho \omega$ in the $c$-space. Its elements therefore have the weight $\omega$. For the iso-parametric motion ( $a=$ const.) the $c$-space is static i.e. each point is fixed. The integral (11) gives us the numerical mean looked for, namely, if $f$ is a phase-function, we have:

$$
\begin{equation*}
F=\int \ldots \int \frac{\partial\left(p_{1} \ldots p_{k}\right)}{\partial\left(c_{1} \ldots c_{k}\right)} d p_{k+1} \ldots d q_{n} / \omega \tag{11'}
\end{equation*}
$$

Returning to equation (7) we now have:

$$
\begin{equation*}
\frac{\partial \varrho}{\partial a}+\sum_{i=1}^{k} \frac{\partial \varrho}{\partial c_{i}} \frac{D c_{i}}{D a}=0 \tag{12}
\end{equation*}
$$

since $\rho$ is a function of $c_{1}, \ldots, c k$ only. Similarly the quantities $D c_{i}{ }_{p a}$ only depend on $c_{1}, \ldots, c_{k}$, as is easily seen from (6), if on the right hand side we replace the time-mean by the numerical mean (11'). Therefore $\varrho$ retains its property $\varrho=\varrho\left(c_{1}, \ldots, c_{k}\right)$ when $a$ changes. Equation (12) expresses, that $e$ is a function of those $k$ integrals of the differential equations (6) which only contain the quantities $c_{1}, \ldots, c_{k}$. These integrals are obtained by integrating the set of $k$ differential equations which on the left side contain the quantities $D c_{i} /_{D a}(i=1,2, \ldots, k)$. They are the essential adiabatic invariants, and we have thus proved that $\varrho$ is a function of the essential adiabatic invariants.

Let us further consider the $c$-space. If $a$ varies slowly, the fixed points in it begin to move. Since in this motion the points do not disappear nor new points are formed, the density go must satisfy the equation of continuity, i, e. our fundamental equation. As $\varrho=$ const. is certainly a possible solution, $\omega$ itself must satisfy the equation

$$
\begin{equation*}
\frac{\partial \omega}{\partial a}+\sum_{i=1}^{k} \frac{\partial \omega \overline{c_{i}^{\prime}}}{\partial c_{i}}=0 \tag{18}
\end{equation*}
$$

where

$$
c_{i}^{\prime}=\frac{D c_{i}}{D a a_{-}}
$$

Or with the notation

$$
\frac{D \omega}{D a}=\frac{\partial \omega}{\partial a}+\sum_{i=1}^{k} \frac{\partial \omega}{\partial c_{i}} D_{i} D_{i}
$$

in the equivalent forms

$$
\begin{equation*}
\frac{D \omega}{D a}+\omega \sum_{i=1}^{k} \frac{\partial c_{i}^{\prime}}{\partial c_{i}}=0, \tag{13'}
\end{equation*}
$$

or

$$
\sum_{i=1}^{k}{\frac{\partial c^{\prime}}{\partial c_{i}}}_{i}=-\frac{1}{\omega} \frac{D \omega}{D a}
$$

For the quantity on the left side - the "divergence" - we shall deduce another expression.
The essential adiabatic invariants - $k$ in number - satisfy the equations

$$
\begin{equation*}
\frac{D v_{i}}{D a_{i}}=\frac{\partial v_{i}}{\partial a}+\sum_{i=1}^{k} \frac{\partial v_{i}}{\partial c_{i}} \bar{c}_{i}^{\prime}(i=1,2, \ldots, k) . \tag{1.4}
\end{equation*}
$$

We shall suppose that the quantities $v_{i}$ can be expressed in the quantities $c_{\lambda}(\lambda=1,2 \ldots k)$ or

$$
\begin{equation*}
\frac{\partial\left(v_{1}, \ldots, v_{k}\right)}{\partial\left(c_{1}, \ldots, c_{k}\right)}=r \neq 0 \tag{15}
\end{equation*}
$$

The properties of our system can be equally well described by the quantities $v_{i}$ as by the quantities $c_{i}$. The ( $v_{1} \ldots v_{k}$ ) -space has the advantage over the $c$-space of being static, also with respect to the action of adiabatic influences. Let us now examine the mutual relation of the two spaces.
To this end we shall consider the $D$-derivative of the determinant $\boldsymbol{r}$ :

$$
\begin{equation*}
\frac{D r}{D a}=\sum_{i \mu} V_{i \mu} \frac{D v_{i \mu}}{D a}, \ldots . \tag{16}
\end{equation*}
$$

where $v_{i \mu}=\frac{\partial v_{i}}{\partial c_{\mu}}$ and $V_{i \mu}$ represents the corresponding sub-determinant. We have identically

$$
\frac{D v_{i}}{D a}=\frac{\partial v_{i}}{\partial a}+\sum_{\lambda} \frac{\partial v_{i}}{\partial c_{\lambda}} \overline{c_{\lambda}^{\prime}}=0
$$

hence

$$
\begin{equation*}
\frac{\partial \nu v_{i}}{\partial c_{\mu}} \frac{\partial^{2} v_{i}}{D a}=\frac{\partial^{2} v_{i}}{\partial a \partial c_{\mu}}+\sum_{\lambda} \frac{\partial c_{\lambda} \partial c_{\mu}}{c_{\lambda}^{\prime}}+\sum_{\lambda} \frac{\partial v_{i} \partial c_{\lambda}^{\prime}}{\partial c_{\lambda}} \frac{\partial c_{\mu}}{}=0 \ldots . \tag{a}
\end{equation*}
$$

and on the other hand

$$
\begin{equation*}
\frac{D}{D a} \frac{\partial v_{i}}{\partial c_{\mu}}=\frac{\partial^{2} v_{i}}{\partial a \partial c_{\mu}}+\sum_{\lambda} \frac{\partial^{2} v_{i}}{\partial c_{\mu} \partial c_{\lambda}} c_{\lambda}^{\prime} \tag{b}
\end{equation*}
$$

From (a) and (b) it follows that

$$
\begin{equation*}
\frac{D \partial v_{i}}{D a} \frac{D v_{i \mu}}{\partial c_{\mu}}=-\sum_{\lambda} \frac{\partial v_{i}}{D a} \frac{\partial c_{\lambda}^{\prime}}{\partial c_{\lambda}} \frac{c_{\mu}}{D a} \tag{17}
\end{equation*}
$$

or

$$
\begin{equation*}
\frac{D v_{i \mu}}{D a}=-\sum_{\lambda=1}^{k} v_{i \lambda} \frac{\partial \overline{c_{\lambda}^{\prime}}}{\partial c_{\mu}} \tag{17}
\end{equation*}
$$

Substituting in (16) we obtain

$$
\frac{D T}{D a}=-\leq v_{i \lambda} V_{i, \mu} \frac{\partial \overline{c_{\lambda}^{\prime}}}{\partial c_{\mu}}
$$

Since $\Sigma v_{i \lambda} V_{i \mu}=0$ for $\lambda \neq \mu$ and equal to $r$ for $\lambda=\mu$, (16) becomes

$$
\frac{D r}{D a}=-r \Sigma \frac{\partial c_{\lambda}^{\prime}}{\partial c_{\lambda}}
$$

Hence

$$
\sum_{\lambda} \sum_{\partial c_{\lambda}^{\prime}}^{\partial c_{\lambda}}=-\frac{1}{r} \frac{D r}{D a}
$$

Comparing this result with (13") it follows that

$$
\begin{equation*}
\frac{1}{\gamma} \frac{D \gamma}{D a}=\frac{1 D \omega}{\omega} . \tag{18}
\end{equation*}
$$

or

$$
\frac{D}{D a} \log \frac{r}{\omega}=0
$$

In other words: $\boldsymbol{\Gamma} / \omega$ is an adiabatic invariant or

$$
\begin{equation*}
r=\omega f\left(v_{1}, \ldots, v_{k}\right), \quad \omega=\boldsymbol{r} F\left(v_{1}, \ldots, v_{k}\right) \tag{18"}
\end{equation*}
$$

Substituting this value of $\omega$ in the integral ( 10 ) we find

$$
\begin{equation*}
I=\int \ldots \int d c_{1} . . d c_{k} r F\left(v_{1}, \ldots, v_{k}\right)=\int \ldots \int d c_{1}, . ., d c_{k} \frac{\partial\left(v_{1} \ldots v_{k}\right)}{\partial\left(c_{1} \ldots c_{k}\right)} F \tag{19}
\end{equation*}
$$

or

$$
I=\int \ldots \int d v_{1} \ldots d v_{k} F\left(v_{1}, \ldots, v_{k}\right)
$$

Now we can always arrange, that $F$ becomes equal to 1 . We have only to introduce, instead of one of the $v_{i}$, the adiabatic invariant

$$
\mathrm{r}_{1}=\text { funct. }\left(v_{1}, \ldots, v_{k}\right)
$$

and submit it to the condition

$$
\begin{equation*}
\frac{\partial r_{1}}{\partial v_{1}}=F\left(v_{1}, \ldots, v_{k}\right) \tag{20}
\end{equation*}
$$

We then find

$$
\begin{gathered}
\frac{\partial r_{1}}{\partial c_{\lambda}}=\sum_{x} \frac{\partial r_{1}}{\partial v_{x}} \frac{\partial v_{x}}{\partial c_{\lambda}}=\sum_{x} \frac{\partial r_{1}}{\partial v_{x}} v_{x \lambda} \\
\boldsymbol{X}^{*}=\frac{\partial\left(r_{1}, v_{2}, \ldots, v_{\gamma k}\right)}{\partial\left(c_{1}, c_{2}, \ldots, c_{k}\right)}=\sum_{\lambda} \frac{\partial r_{1}}{\partial c_{\lambda}} V_{1 \lambda}=\Sigma \frac{\partial r_{1}}{\partial v_{x}} V_{i \lambda} v_{x \lambda}=\frac{\partial r_{1}}{\partial v_{1}} \boldsymbol{r}
\end{gathered}
$$

or substituting for $x$ its value $\omega / F$ :

$$
x^{*}=\omega \frac{\partial r_{1}}{\partial v_{1}} \frac{F}{1}=\omega
$$

Calling the thus normalized set of essential adiabatic invariants $v_{1}, \ldots, v_{k}$, we find

$$
\begin{equation*}
I=\int \ldots \int d v_{1} \ldots d v_{k} \tag{22}
\end{equation*}
$$

The $v$-space which is static with respect to adiabatic action is "weightless": its density $\rho$ is simply equal to $\varrho\left(v_{1}, \ldots, v_{k}\right)$. The quantities $v_{1}, \ldots, v_{k}$ may be quanticized. Its property which is expressed by eq. (22) is nothing but the fundamental law which according to Pranck's hypothesis the quanta-quantities have to obey ${ }^{1}$ ). By our theorem ( $18^{\prime \prime}$ ) this hypothesis is connected with the adiabatic invariants and thus finds a new confirmation. The property of the $v$-space being "weightless" displays the character of this fundamental law as a natural generalisation of the old quanta-hypothesis.
§4. On the coherence of degrees of freedom.
From the point of view now attained this very important conception appears as a natural consequence of our suppositions. If the number $k$-. that of the essential integrals and adiabatic invariants - is smaller than the number of degrees of freedom $n$, as appears from (22), some of the quantities $v_{i}$ must necessarily be of the dimen-

[^3]sion $h^{p}(p>1)$, since the dimension of $I$ is $h^{n}$. In order to illustrate this and the previous results we shall contrast the properties of a Boitzmann ergode and a conditionally periodic system without commensurable relations:
\[

$$
\begin{aligned}
& \text { ergode . conditional periodic system } \\
& \text { essential integrals } \\
& H=c_{1} \quad \because \quad H_{1}=c_{1}, H_{3}=c_{2}, \ldots, H_{n}=c_{n} . \\
& \text { numerical mean } \\
& \int \ldots \int d p_{2} \ldots d p_{n} d q_{1} \ldots d q_{n} \frac{\partial p_{1}}{\partial c_{1}} \cdot(23) \left\lvert\, \underset{\text { density }}{\int} \ldots \int d q_{1} \ldots d q_{n} \frac{\partial\left(p_{1}, \ldots, p_{n}\right)}{\partial\left(c_{1}, \ldots, c_{n}\right)}\left(23^{\prime}\right)\right. \\
& \varrho=\varrho\left(c_{1}, a\right) \quad \varrho \equiv \varrho\left(c_{1}, \ldots, c_{n} ; a\right)
\end{aligned}
$$
\]

essential adiabatic invariants

$$
\left.\begin{gathered}
V=\int_{\frac{H}{H \leq c_{1}}}^{\int} \ldots \int_{1} d p_{1} \ldots d q_{n}
\end{gathered} \quad\right|_{i}=\int_{0}^{\infty} p_{i} d q_{i}(i=1,2, \ldots, n) .
$$

The conditionally periodic system is what Bolizmann calls a subergode. On the other hand the ergode appears as a coherent system with a smallest value of $k$, viz. $k=1$. These short indications may suffice for the present.

## § 5. Degeneration.

A conditionally periodic system is called degenerated, if there are commensurable relations between the periodicity moduli. It is evident, that our system covers a lower set of points with its orbital curve everywhere densely, than when there are no such relations. Accordingly the numerical mean will be of a lower dimension and more quantities will remain free after the averaging process. Thus besides the quantities $c$ the quantities $t$ will play a part: the number of essential adiabatic invariants becomes larger than the number of degrees of freedom. The question, whether these supernumerary quantities have to be quanticized, we shall not discuss here. A good instance for the discussion of the questions which may arise here is afforded by the quanta-quantities in Epstein's theory ${ }^{\text {i }}$ ) of the Stark-effect for an infinitely weak external electric field; the "parabolic" quanta-quantities which are found in this case cannot

[^4]be represented as function of Sommerfeld's "spherical" quanta-quantities alone; other adiabatic invariants containing the quantities $t_{2}$ and $t_{\mathrm{t}}$ are essential in this case.
I am fully conscious of the fact, that by the above considerations the difficulties which still beset the theory of quanta are in no way removed, but only shifted. Still it seems to me that even the possibility of such displacement deserves attention. Moreover I expect that in special cases the general theory tentatively sketched out here may be found useful.

Physical Laboratory of the University. Petrograd, April 1, 1920.


[^0]:    ${ }_{2}^{1)}$ M. Planck. Ann. d. Phys. 50 (1916) p. 392.
    ${ }^{\text {2 }}$ ) M. Planck, l.c.
    ${ }^{\text {8 }}$ ) K. Schwarzschild. Sitzungsber. Berlin 1916. P. Epstein Ann. d. Phys, 51. (1916). p. 168.

[^1]:    ${ }^{1}$ ) In the theory of the Zeemaneeffect as given by Sommerfeld and Debije (Phys. Zschr. 17. (1916) a difficulty is met with here. This may, I think, be evaded in different ways, but I am not able to give a uniquely determined solution.

[^2]:    ${ }^{1}$ ) For the method now following comp. J. W. Gibbs. Statistical Mechanics Chapter XIII.
    ${ }^{\text {2 }}$ ) Comp. J. M. Burgers. Proc. Amst. XX (1916) 149, Ann. d. Phys. 1917 (2) and my paper in the Proc. Amst. 1. c.

[^3]:    ${ }^{1}$ ) M. P PaNCK. .f.

[^4]:    ${ }^{1}$ ) P. Epstein. Ann. d. Phys. 50. p. 490.

